On stability of difference schemes. Central schemes for hyperbolic conservation laws with source terms

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Abstract

The stability of nonlinear explicit difference schemes with not, in general, open domains of the scheme operators are studied. For the case of path-connected, bounded, and Lipschitz domains, we establish the notion that a multi-level nonlinear explicit scheme is stable *iff* (if and only if) the corresponding scheme in variations is stable. A new modification of the central Lax-Friedrichs (LxF) scheme is developed to be of the second order accuracy. The modified scheme is based on nonstaggered grids. A monotone piecewise cubic interpolation is used in the central scheme to give an accurate approximation for the model in question. The stability of the modified scheme is investigated. Some versions of the modified scheme are tested on several conservation laws, and the scheme is found to be accurate and robust. As applied to hyperbolic conservation laws with, in general, stiff source terms, it is constructed a second order nonstaggered central scheme based on operator-splitting techniques.

1 Introduction

We are mainly concerned with the stability of nonlinear explicit difference schemes arising, e.g., in numerical analysis of nonlinear PDE (partial differential equation, such as parabolic, hyperbolic, etc.) systems, generally, in multidimensional space. Stability is the central and the most pressing problem in any algorithm [48]. Nowadays, there exists a few methods for stability analysis of some classes of difference schemes approximating PDE systems (see, e.g., [16], [17], [32], [35], [38], [48] and references therein). However, the problem of stability analysis for nonlinear schemes is still one of the most burning problems,

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because of the absence of its complete solution [17]. In particular, the vast majority of difference schemes, currently in use, have still not been analyzed [16], and, in general, no numerical method for non-linear systems of equations has been proven to be stable [32]. A new approach to testing scheme stability has been suggested in [7] to prove convergence of non-linear schemes for systems of PDEs. It was demonstrated [7] that the notion of scheme in variations (or variational scheme [9], [10]) has much potential to be an effective tool for studying stability of nonlinear schemes. The scheme in variations for a difference scheme represents a tangent space at a point of the manifold associated to the difference scheme [7], [9]. Thus, the scheme in variations will always be linear and, hence, enables the investigation of the stability for nonlinear operators using, in general, an infinite family of linear patterns. It was established in [7] the notion that the stability of a scheme in variations implies the stability of its original scheme, and that a nonlinear explicit scheme will be stable iff (if and only if) its scheme in variations will be stable. In this paper we consider difference schemes with not, in general, open domains of the scheme operators. For the case of internal path-connected, bounded, and Lipschitz domains, we establish the notion that a multi-level nonlinear explicit scheme is stable iff the corresponding scheme in variations is stable.

By way of illustration of the developed approach we are concerned with the stability analysis of central difference schemes for hyperbolic systems of conservation laws with source terms. Such systems are used to describe many physical problems of great practical importance in magneto-hydrodynamics, kinetic theory of rarefied gases, linear and nonlinear waves, viscoelasticity, multi-phase flows and phase transitions, shallow waters, etc. (see, e.g., [6], [12], [18], [24], [29], [32], [34], [37], [41], [42]). We will consider a system of hyperbolic conservation laws written as follows (e.g., [18], [32])

$$\frac{\partial \mathbf{u}}{\partial t} + \sum_{j=1}^{N} \frac{\partial}{\partial x_{j}} \mathbf{f}_{j} (\mathbf{u}) = \frac{1}{\tau} \mathbf{q} (\mathbf{u}), \ 0 < t \le T_{\text{max}}, \ \mathbf{u} (\mathbf{x}, t)|_{t=0} = \mathbf{u}^{0} (\mathbf{x}), \quad (1)$$

where the column-vector $\mathbf{x} \equiv \{x_1, x_2, \dots, x_N\}^T \in \mathbb{R}^N$, $\mathbf{u} = \{u_1, u_2, \dots, u_M\}^T$ is a vector-valued function from $\mathbb{R}^N \times [0, +\infty)$ into a subset $\Omega_{\mathbf{u}} \subset \mathbb{R}^M$, $\mathbf{f}_j(\mathbf{u}) = \{f_{1j}(\mathbf{u}), f_{2j}(\mathbf{u}), \dots, f_{Mj}(\mathbf{u})\}^T$ is a smooth function (flux-function) from $\Omega_{\mathbf{u}}$ into \mathbb{R}^M , $\mathbf{q}(\mathbf{u}) = \{q_1(\mathbf{u}), q_2(\mathbf{u}), \dots, q_M(\mathbf{u})\}^T$ denotes the source term, $\tau > 0$ denotes the stiffness parameter, $\mathbf{u}^0(\mathbf{x})$ is of compact support. It is assumed that $\tau = const$ without loss of generality. It is also assumed that all eigenvalues, $\xi_k = \xi_k(\mathbf{u})$, of the Jacobian matrix $\mathbf{G}(=\partial \mathbf{q}(\mathbf{u})/\partial \mathbf{u})$ have non-positive real parts, i.e.

$$R_e \xi_k \left(\mathbf{u} \right) \le 0, \quad \forall k, \quad \forall \mathbf{u} \in \Omega_{\mathbf{u}}.$$
 (2)

In what follows $\|\mathbf{M}\|_p$ denotes the matrix norm of a matrix \mathbf{M} induced by the vector norm $\|\mathbf{v}\|_p = (\sum_i |v_i|^p)^{1/p}$, and $\|\mathbf{M}\|$ denotes the matrix norm induced by a prescribed vector norm. \mathbb{R} denotes the field of real numbers. The transpose of a matrix (or vector) \mathbf{S} is denoted by \mathbf{S}^T . The null element in any linear space, as well as the number zero, will be denoted by the same symbol 0.

Central schemes are attractive for various reasons: no Riemann solvers, characteristic decompositions, complicated flux splittings, etc., must be involved in construction of a central scheme (see, e.g., [5], [8], [23], [29], [30], [32], [39], [41], [43] and references therein). However, there is the need for staggered central schemes to alternate between two staggered grids, which can be cumbersome, e.g., near the boundaries [23]. Moreover, there is a risk for every central scheme to exhibit spurious solutions [8] in spite of sufficiently small CFL (Courant-Friedrichs-Lewy [32]) number. To simplify implementation of central schemes in the case of complex geometries and boundary conditions, it was developed nonstaggered central schemes in [23]. A new second-order modification of the staggered LxF scheme was developed in [8] to avoid the risk of spurious oscillations. It was demonstrated in [8] that the higher order versions of LxF scheme can produce spurious oscillations because of a negative numerical viscosity introduced to increase accuracy of the scheme up to $O(\Delta t + (\Delta x)^2)$. To reduce the risk of spurious solutions, it was suggested in [8] to introduce an additional non-negative numerical viscosity such that the scheme's order of accuracy becomes $O((\Delta t)^2 + (\Delta x)^2)$. The developed central scheme was tested on several conservation laws taking a CFL number equal or close to unity, and the scheme was found to be accurate and robust. In this paper, we extend the second-order modification to a new nonstaggered central scheme. The stability of this scheme is proven in Section 3.3 on the basis of the approach developed in Section 2 as well as in [7], [8]. The scheme is tested on several conservation laws in Section

A stable numerical scheme may yield spurious results when applied to a stiff hyperbolic system with relaxation (see, e.g., [1], [4], [6], [12], [13], [24], [44], [45]). It is significant that a numerical scheme for relaxation systems must possess a discrete analogy to the continuous asymptotic limit, because any scheme violating the correct asymptotic limit leads to spurious or poor solutions (see, e.g., [12], [24], [25], [37], [41]). Most methods for solving such systems can be described as operator splitting ones, [13], or methods of fractional steps, [6]. After operator splitting, one solves the advection homogeneous system, and then the ordinary differential equations associated with the source terms. We are mainly concerned with such an approach in Section 4.

2 Stability of nonlinear explicit schemes

We consider the following (q + 1)-level difference scheme:

$$\mathbf{w}_{i}^{n+1} = \mathbf{G}_{i}^{n}(\mathbf{w}_{1}^{n}, \mathbf{w}_{1}^{n-1}, \dots, \mathbf{w}_{1}^{n-q+1}, \mathbf{w}_{2}^{n}, \mathbf{w}_{2}^{n-1}, \dots, \mathbf{w}_{2}^{n-q+1}, \dots,$$

$$\mathbf{w}_{I}^{n}, \mathbf{w}_{I}^{n-1}, \dots, \mathbf{w}_{I}^{n-q+1}), \quad \mathbf{G}_{i}^{n} : \Omega_{n} \subseteq \mathbb{R}^{N} \to \mathbb{R}^{N_{0}},$$

$$(3)$$

where $q, N_0, I \geq 1$ are finite integer constants, $N = qN_0I$, $\mathbf{w}_i^k \in \mathbb{R}^{N_0}$ denotes a vector-valued grid function, $i \in \omega_1$ denotes a node of the grid $\omega_1 \equiv \{1, 2, \dots, I\}$, $k \in \omega_2$ denotes a node (time level) of the grid $\omega_2 \equiv \{0, 1, \dots, M\}$, $M \geq q$ is a finite integer constant. $\mathbf{G}_i^n = \left\{G_{i,1}^n, \ G_{i,2}^n, \dots, \ G_{i,N_0}^n\right\}^T$ denotes the vector-valued

function with the domain and range belonging to \mathbb{R}^N and \mathbb{R}^{N_0} , respectively. We emphasize that the domain, Ω_n , of the function \mathbf{G}_i^n is the set of argument values, which need not be open in \mathbb{R}^N . We assume that n in (3) denotes the time level, t_n (= $n\Delta t$). Thus, the time increment will be represented by $\Delta t = t_{\max}/M = const$, where t_{\max} denotes some finite time over which we wish to compute. Using the notation

$$\mathbf{H}_{i}^{n} = \left\{ \left(\mathbf{G}_{i}^{n}\right)^{T}, \left(\mathbf{w}_{i}^{n}\right)^{T}, \left(\mathbf{w}_{i}^{n-1}\right)^{T}, \dots, \left(\mathbf{w}_{i}^{n-q+2}\right)^{T} \right\}^{T},$$

$$\mathbf{v}_{i}^{n} = \left\{ \left(\mathbf{w}_{i}^{n}\right)^{T}, \left(\mathbf{w}_{i}^{n-1}\right)^{T}, \dots, \left(\mathbf{w}_{i}^{n-q+1}\right)^{T} \right\}^{T}, \quad i \in \omega_{1}, \ n \in \omega_{2},$$

$$(4)$$

we reduce (see, e.g., [15], [47]) the multi-level scheme, (3), to the following two-level scheme:

$$\mathbf{v}_i^{n+1} = \mathbf{H}_i^n(\mathbf{v}_1^n, \mathbf{v}_2^n, \dots, \mathbf{v}_I^n), \ \mathbf{H}_i^n : \Omega_n \subseteq \mathbb{R}^N \to \mathbb{R}^{qN_0}, \ i \in \omega_1, \ n, n+1 \in \omega_2.$$
 (5)

If we define

$$\mathbf{v}^{n} = \left\{ (\mathbf{v}_{1}^{n})^{T}, (\mathbf{v}_{2}^{n})^{T}, \dots, (\mathbf{v}_{I}^{n})^{T} \right\}^{T}, \ \mathbf{H}^{n} = \left\{ (\mathbf{H}_{1}^{n})^{T}, (\mathbf{H}_{2}^{n})^{T}, \dots, (\mathbf{H}_{I}^{n})^{T} \right\}^{T}, \ (6)$$

then the scheme (5) becomes

$$\mathbf{v}^{n+1} = \mathbf{H}^n(\mathbf{v}^n), \quad \mathbf{H}^n : \Omega_n \subseteq \mathbb{R}^N \to \mathbb{R}^N, \quad n, n+1 \in \omega_2 \equiv \{0, 1, \dots, M\}.$$
 (7)

As usual (e.g., [40, p. 62]), for mappings $\mathbf{f}: \Omega_f \subseteq \mathbb{R}^N \to \mathbb{R}^N$ and $\mathbf{g}: \Omega_g \subseteq \mathbb{R}^N \to \mathbb{R}^N$, the composite mapping $\mathbf{h} = \mathbf{g} \circ \mathbf{f}$ is defined by $\mathbf{h}(\mathbf{v}) = \mathbf{g}(\mathbf{f}(\mathbf{v}))$ for all $\mathbf{v} \in \Omega_h = {\mathbf{v} \in \Omega_f \mid \mathbf{f}(\mathbf{v}) \in \Omega_g}$. Using the composite mapping approach, we rewrite Scheme (7) to read

$$\mathbf{y} = \mathbf{F}(\mathbf{x}), \quad \mathbf{F} : \Omega_F \subseteq \mathbb{R}^N \to \mathbb{R}^N,$$
 (8)

where the following notation is used: $\mathbf{x} = \mathbf{v}^0$, $\mathbf{y} = \mathbf{v}^M$, $\mathbf{F} = \mathbf{H}^{M-1} \circ \mathbf{H}^{M-2} \circ \ldots \circ \mathbf{H}^0$, $\Omega_F = \{\mathbf{v}^0 \in \Omega_0 \mid \mathbf{v}^1 = \mathbf{H}^0 (\mathbf{v}^0) \in \Omega_1 \mid \ldots \mid \mathbf{v}^{M-1} = \mathbf{H}^{M-2} (\mathbf{v}^{M-2}) \in \Omega_{M-1} \}$. Notice, \mathbf{F} in (8) depends also on scheme parameters (e.g., space and time increments), however, this dependence is usually not included in the notation. Let the scheme parameters (including time increments) be represented by a vector \mathbf{s} belonging to some normed space with the norm $|\mathbf{s}|$.

Let us remind the well-known definition of stability of an explicit scheme (see, e.g., [16], [18], [32], [47], [48], [49]).

Definition 1 Scheme (8) is said to be stable if there exist positive constants s_0 and C such that for all \mathbf{x} , $\mathbf{x}_* \in \Omega_F$ the following inequality is valid

$$\|\mathbf{F}(\mathbf{x}_*) - \mathbf{F}(\mathbf{x})\| \le C \|\mathbf{x}_* - \mathbf{x}\|, \quad \forall \ \mathbf{s} : \ |\mathbf{s}| \le s_0.$$
 (9)

Thus, scheme (8) and, hence, scheme (3) will be stable *iff* the function **F** will be Lipschitz for a constant C. Notice that this well-known definition of stability may lead to wrong conclusion. Actually, let us consider the "slit plane" [21] in polar coordinates (r, θ)

$$\Omega_F = \{ (r, \theta) \mid 0 < r < \infty, -\pi < \theta < \pi \} \subset \mathbb{R}^2, \tag{10}$$

and the vector-valued function $\mathbf{F} = \{F_1(r,\theta), F_2(r,\theta)\}^T$ such that [21]

$$F_1 = r, \ F_2 = \theta/2.$$
 (11)

We take $(r,\theta)_* = (r_0, -\pi + \varepsilon)$, $(r,\theta) = (r_0, \pi - \varepsilon)$, and $r_0 = const$. Considering the Euclidean distance between these two points of the plane with Cartesian coordinates $\mathbf{x}_* \equiv (-r_0 \cos \varepsilon, -r_0 \sin \varepsilon)$ and $\mathbf{x} \equiv (-r_0 \cos \varepsilon, r_0 \sin \varepsilon)$ we find that the inequality (9) takes the form: $\cos(\varepsilon/2) \leq C \sin \varepsilon$. Obviously the mapping (11) is not Lipschitz, since C in (9) tends to infinity as $\varepsilon \to 0$. Therefore, we have to conclude, in view of the above definition, that the scheme (11) is not stable, even though for every point (r,θ) in Ω_F , (10), there exists a neighborhood of (r,θ) such that the function \mathbf{F} , (11), restricted to the neighborhood is Lipschitz for C=1, i.e. the function \mathbf{F} is locally Lipschitz. Further, \mathbf{F} is the non-stretching mapping of the "slit plane" (10) into the right semi-plane. Hence, the preceding definition of stability, i.e. Definition 1, needs to be slightly improved. Such a problem has been considered in [7] under the assumption that the domain Ω_F of the vector-valued function \mathbf{F} , (8), is open in \mathbb{R}^N . We will assume that the Lebesgue-measurable domain Ω_F of the function \mathbf{F} , (8), is a not, in general, open subset of \mathbb{R}^N with non-empty interior.

We shall start with some basic notions. A set $\Omega \subseteq \mathbb{R}^N$ is said to be internal path-connected (or path-connected, in the case that Ω is open) if every two points $\mathbf{x}, \mathbf{x}_* \in \Omega$ can be joined by a continuous curve $\gamma : [0,1] \subset \mathbb{R} \to \Omega$ of finite length, $L(\gamma)$, with $int(\gamma) \subseteq int(\Omega)$. Here and in what follows $int(\Omega)$ denotes the interior of $\Omega \subseteq \mathbb{R}^N$ and $int(\gamma) \equiv \gamma \setminus (\gamma(0) \cup \gamma(1))$ denotes the interior of the continuous curve γ . The intrinsic metric Λ_{Ω} on an internal path-connected set Ω is defined as

$$\Lambda_{\Omega}(\mathbf{x}, \mathbf{x}_{*}) = \inf_{\gamma \subseteq \Omega} L(\gamma), \quad \mathbf{x} = \gamma(0), \ \mathbf{x}_{*} = \gamma(1), \ \mathbf{x}, \mathbf{x}_{*} \in \Omega.$$
 (12)

An open ball (of radius r > 0) about $\mathbf{x} \in \mathbb{R}^N$ is denoted by $B(\mathbf{x}, r)$ (or just $B_{\mathbf{x}}$). A closed ball is denoted by $\overline{B}(\mathbf{x}, r)$ (or just $\overline{B}_{\mathbf{x}}$). Given points $\mathbf{x}, \mathbf{y} \in \mathbb{R}^N$, given an open ball $B(\mathbf{x}, r_{\mathbf{x}})$, and given an open ball $B(\mathbf{y}, r_{\mathbf{y}})$ such that $\mathbf{x} \notin B_{\mathbf{y}}$, the convex open set $Q_{\mathbf{x}} = \{\mathbf{z} \mid \mathbf{x} + \lambda(\mathbf{z} - \mathbf{x}), \lambda > 0, \mathbf{z} \in B_{\mathbf{y}}\} \cap B_{\mathbf{x}}$ will be called a finite cone with vertex \mathbf{x} . A set $\Omega \subseteq \mathbb{R}^N$ is said to have the cone property if there exists a finite cone $Q_{\mathbf{x}}$ such that each boundary point $\mathbf{x} \in \partial \Omega$ is the vertex of a finite cone $Q_{\mathbf{x}} \subseteq int(\Omega)$ congruent to $Q_{\mathbf{x}}$. Suppose that $\Omega \subseteq \mathbb{R}^N$ has the cone property, and let $K_{\mathbf{x}}$ denote the union of all finite cones with the vertex $\mathbf{x} \in \partial \Omega$, i.e. $K_{\mathbf{x}} = \cup Q_{\mathbf{x}}$. Then a function $\mathbf{F} : \Omega \subseteq \mathbb{R}^N \to \mathbb{R}^N$ will be called locally Lipschitz (or locally Lipschitz continuous) if for every $\mathbf{x} \in \Omega$ there exist $B_{\mathbf{x}}$ and a constant $c_{\mathbf{x}} > 0$ such that $\|\mathbf{F}(\mathbf{x}) - \mathbf{F}(\mathbf{y})\| \le c_{\mathbf{x}} \|\mathbf{x} - \mathbf{y}\|$ for each

 $\mathbf{y} \in B_{\mathbf{x}}$ provided $\mathbf{x} \in int(\Omega)$, or $\mathbf{y} \in B_{\mathbf{x}} \cap K_{\mathbf{x}}$ otherwise. $\Omega \subset \mathbb{R}^N$ is Lipschitz (has a Lipschitz boundary), if for each point $\mathbf{x} \in \partial \Omega$ there exists an open ball $B_{\mathbf{x}}$ such that $\partial \Omega \cap B_{\mathbf{x}}$ is the graph of a Lipschitz function and $int(\Omega)$ lies on one side of its boundary (for more details see, e.g., [11, p. 149], [31, p. 354]). Let us note that $\Omega \subset \mathbb{R}^N$ satisfies the cone property if it has a Lipschitz boundary (see, e.g., [11, p. 151], [31, p. 355]), but not vice-versa as the example (10) shows. $W^{1,\infty}(\Omega)$ denotes the Sobolev space (see Definition 5 in [11, p. 28], see also [21] and [31]), i.e., the space consisting of all bounded functions on open $\Omega \subset \mathbb{R}^N$ whose distributional derivatives are bounded functions as well. Let $\mathbf{F} \equiv \{F_1, \dots, F_N\}^T$ in (8), let $\nabla F_i \equiv \{\partial_1 F_i, \dots, \partial_N F_i\}$ denote the distributional gradient of F_i , and let $\delta \mathbf{F}, \delta \mathbf{x} \in \mathbb{R}^N$ denote variations. The following equality

$$\delta \mathbf{F} = \mathbf{F}' \cdot \delta \mathbf{x}, \quad \mathbf{F}' \equiv \left\{ (\nabla F_1)^T, (\nabla F_2)^T, \dots, (\nabla F_N)^T \right\}^T,$$
 (13)

will be viewed as the scheme in variations for (8). The matrix representation of \mathbf{F}' in (13) is given by the following Jacobian matrix [40]: $\mathbf{F}' = \{\partial F_i / \partial x_j\}, i, j = 1, 2, ..., N$.

Definition 2 Let Ω_F in (8) be internal path-connected. Scheme (8) is said to be stable if there exist positive constants s_0 and C such that the following inequality holds

$$\|\mathbf{F}(\mathbf{x}_*) - \mathbf{F}(\mathbf{x})\| \le C\Lambda_{\Omega_F}(\mathbf{x}, \mathbf{x}_*), \quad \forall \ \mathbf{x}, \mathbf{x}_* \in \Omega_F, \quad \forall \ \mathbf{s} : \ |\mathbf{s}| \le s_0.$$
 (14)

Notice that the scheme (11) is stable in terms of Definition 2, since the inequality (14) holds for C=1. It is significant that the domain Ω_F , (10), of the function \mathbf{F} , (11), is open in \mathbb{R}^2 . However, if the domain of the function \mathbf{F} , (11), is not open, then the situation can be in complete contrast to the previous one. Actually, let us add to Ω_F , (10), the only boundary point $(1,\pi)$, i.e. let $\Omega_F^* = \Omega_F \cup (1,\pi)$ will be the domain of the function $\mathbf{F}^* = \mathbf{F}$, (11), if $(r,\theta) \in \Omega_F$ and $\mathbf{F}^* = (1,\pi/2)$ if $(r,\theta) = (1,\pi)$. We take $(r,\theta)_* = (1,-\pi+\varepsilon)$, $(r,\theta) = (1,\pi)$. Considering the Euclidean distance between these two points of the plane with Cartesian coordinates $\mathbf{x}_* \equiv (-\cos\varepsilon, -r_0\sin\varepsilon)$ and $\mathbf{x} \equiv (-1,0)$ we find that the inequality (14) takes the form: $\cos(\varepsilon/4) \leq C\sin(\varepsilon/2)$. Thus, the scheme (11), namely $\mathbf{F} : \Omega_F^* \subset \mathbb{R}^2 \to \mathbb{R}^2$, is not stable, since C in (14) tends to infinity as $\varepsilon \to 0$.

Lemma 3 Let the path-connected Ω_F of (8) be open in \mathbb{R}^N . Scheme (8) will be stable in terms of Definition 2 iff \mathbf{F} in (8) will be locally Lipschitz for a common constant C, for all scheme parameters \mathbf{s} such that $|\mathbf{s}| \leq s_0$.

Proof. Suppose Scheme (8) is stable, i.e. (14) is valid. Choose any point $\mathbf{x} \in \Omega_F$. Since Ω_F is open, there exists a radius r such that $B(\mathbf{x}, r) \subset \Omega_F$. Choose any point $\mathbf{x}_* \in B(\mathbf{x}, r)$, and let γ_* be the straight line segment joining the points $\mathbf{x}, \mathbf{x}_* \in B(\mathbf{x}, r)$. In view of (14), \mathbf{F} in (8) will be locally Lipschitz for a common constant C, for all $\mathbf{s} : |\mathbf{s}| \leq s_0$, since $\Lambda_{\Omega_F}(\mathbf{x}, \mathbf{x}_*) = L(\gamma_*) = ||\mathbf{x}_* - \mathbf{x}||$.

Conversely, suppose that \mathbf{F} in (8) is locally Lipschitz for a common constant C, for all $\mathbf{s}: |\mathbf{s}| \leq s_0$. Let some points $\mathbf{x}, \mathbf{x}_* \in \Omega_F$ be joined by a continuous curve γ . In view of (12), the curve γ can be taken such that $L(\gamma) \leq \Lambda_{\Omega_F}(\mathbf{x}, \mathbf{x}_*) + \varepsilon$ for an arbitrary $\varepsilon > 0$. Given any point $\mathbf{z} \in \gamma$, there is a ball $B_{\mathbf{z}} \subset \Omega_F$ such that \mathbf{F} restricted to $B_{\mathbf{z}}$ is Lipschitz for the common constant C (locally Lipschitz continuity). The balls $\{B_{\mathbf{z}}\}$ form an open cover of γ . Since the mapping $\gamma:[0,1] \subset \mathbb{R} \to \mathbb{R}^N$ is continuous, the curve γ is compact [28, p. 94]. Hence, by the compactness of γ , $\{B_{\mathbf{z}}\}$ has a finite subcover consisting of balls $B_{\mathbf{x}} = B_{\mathbf{z}_1}, B_{\mathbf{z}_2}, \ldots, B_{\mathbf{z}_K} = B_{\mathbf{x}_*}$. Since \mathbf{F} is locally Lipschitz, we find

$$\|\mathbf{F}(\mathbf{z}_{k+1}) - \mathbf{F}(\mathbf{z}_k)\| \le C \|\mathbf{z}_{k+1} - \mathbf{z}_k\|, \quad k = 1, 2, \dots K - 1, \ \forall \ \mathbf{s} : \ |\mathbf{s}| \le s_0.$$
 (15)

Then, by virtue of (15), we find

$$\|\mathbf{F}(\mathbf{x}_{*}) - \mathbf{F}(\mathbf{x})\| = \left\| \sum_{k} \left[\mathbf{F}(\mathbf{z}_{k+1}) - \mathbf{F}(\mathbf{z}_{k}) \right] \right\| \le C \sum_{k} \|\mathbf{z}_{k+1} - \mathbf{z}_{k}\| \le CL(\gamma) \le C\Lambda_{\Omega_{F}}(\mathbf{x}, \mathbf{x}_{*}) + \varepsilon C, \quad \forall \ \mathbf{s} : \ |\mathbf{s}| \le s_{0}.$$
 (16)

By letting $\varepsilon \to 0$, we find that (14) holds. \blacksquare

Notice, in view of Lemma 3, the scheme (11) is stable in terms of Definition 2, since the domain Ω_F , (10), is open in \mathbb{R}^2 and path-connected, and the function \mathbf{F} , (11), is locally Lipschitz for the common constant C = 1.

Theorem 4 Let the internal path-connected Ω_F of (8) have the cone property, then the scheme (8) will be stable in terms of Definition 2 iff \mathbf{F} in (8) will be locally Lipschitz for a common constant C, for all \mathbf{s} such that $|\mathbf{s}| \leq s_0$.

Proof. Let $K_{\mathbf{x}}$ denote the union of all finite cones with the vertex $\mathbf{x} \in \partial \Omega_F$. Suppose Scheme (8) is stable. Choose any point $\mathbf{x} \in \Omega_F$. If $\mathbf{x} \in int(\Omega_F)$, then, in view of Lemma 3, \mathbf{F} in (8) will be locally Lipschitz. If $\mathbf{x} \notin int(\Omega_H)$, then (cone property) for each $\mathbf{x}_* \in K_{\mathbf{x}}$ there is a straight line segment joining the points \mathbf{x} , \mathbf{x}_* . Since $\Lambda_{\Omega_F}(\mathbf{x}, \mathbf{x}_*) = \|\mathbf{x}_* - \mathbf{x}\|$, the necessity, in view of (14), is proven.

Conversely, suppose \mathbf{F} in (8) is locally Lipschitz. Let some points \mathbf{x} , $\mathbf{x}_* \in \Omega_F$ be joined by a continuous curve γ . If \mathbf{x} , $\mathbf{x}_* \in int(\Omega_F)$, then, in view of Lemma 3, we find that (14) holds. Let $\mathbf{x} \notin int(\Omega_F)$, however $\mathbf{x}_* \in int(\Omega_F)$. Choose any point $\mathbf{z} \in K_{\mathbf{x}}$ such that $\|\mathbf{x} - \mathbf{z}\| = \varepsilon$. Then, since \mathbf{F} is locally Lipschitz, we have $\|\mathbf{F}(\mathbf{x}) - \mathbf{F}(\mathbf{z})\| \le C \|\mathbf{x} - \mathbf{z}\|$ for a sufficiently small ε . Since $\mathbf{z} \in int(\Omega_F)$, we find, in view of Lemma 3, that $\|\mathbf{F}(\mathbf{z}) - \mathbf{F}(\mathbf{x}_*)\| \le C\Lambda_{\Omega_F}(\mathbf{z}, \mathbf{x}_*)$. Then

$$\|\mathbf{F}(\mathbf{x}) - \mathbf{F}(\mathbf{x}_*)\| \le \|\mathbf{F}(\mathbf{x}) - \mathbf{F}(\mathbf{z})\| + \|\mathbf{F}(\mathbf{z}) - \mathbf{F}(\mathbf{x}_*)\|$$

$$\le \varepsilon C + \|\mathbf{F}(\mathbf{z}) - \mathbf{F}(\mathbf{x}_*)\|, \quad \forall \ \mathbf{z} : \ \|\mathbf{x} - \mathbf{z}\| = \varepsilon, \ \forall \ \mathbf{s} : \ |\mathbf{s}| \le s_0. \tag{17}$$

Since $\|\mathbf{F}(\mathbf{x}) - \mathbf{F}(\mathbf{x}_*)\|$ does not depend on \mathbf{z} , we find

$$\|\mathbf{F}(\mathbf{x}) - \mathbf{F}(\mathbf{x}_*)\| \le \inf_{\mathbf{z}: \|\mathbf{x} - \mathbf{z}\| = \varepsilon} \{\varepsilon C + \|\mathbf{F}(\mathbf{z}) - \mathbf{F}(\mathbf{x}_*)\|\}$$
 (18)

$$\leq \varepsilon C + C \inf_{\mathbf{z}: \|\mathbf{x} - \mathbf{z}\| = \varepsilon} \Lambda_{\Omega_F} (\mathbf{z}, \mathbf{x}_*), \quad \forall \ \mathbf{s}: \ |\mathbf{s}| \leq s_0.$$
 (19)

By letting $\varepsilon \to 0$, i.e. $\mathbf{z} \to \mathbf{x}$, we find that (14) holds. The proof in the case \mathbf{x} , $\mathbf{x}_* \notin int(\Omega_F)$ is reduced to the previous one by choosing $\mathbf{z} \in K_{\mathbf{x}_*}$.

We shall need the following theorem that identifies the Sobolev space $W^{1,\infty}(\Omega)$ as a space of locally Lipschitz functions.

Theorem 5 Let $\Omega_F \subseteq \mathbb{R}^N$ be open, and let $\mathbf{F} = \{F_1, F_2, \dots, F_N\}^T$ in (8). Then, F_i , $i = 1, 2, \dots, N$, (and, hence, \mathbf{F}) is locally Lipschitz (in the sense of having representatives) iff $F_i \in W^{1,\infty}(\Omega_F)$.

The proof of Theorem 5 can be found, e.g., in [21, Theorem 4.1], [31, p. 342]. The following lemma gives the necessary and sufficient conditions for the stability of the scheme in variations (13).

Lemma 6 Linear Scheme (13) will be stable iff there exist positive s_0 ,C = const such that

$$\|\mathbf{F}'\| \le C = const, \quad \forall \ \mathbf{x} \in \Omega_F, \quad \forall \ \mathbf{s} : \ |\mathbf{s}| \le s_0.$$
 (20)

The proof of Lemma 6 can be found in [7] (see also, e.g., [16], [47], [48], [49]). Let us remind the well-known theorem that ensures that Lipschitz domains are extension domains (see [31, pp. 320, 356], [11, pp. 149, 285]).

Theorem 7 Let $\Omega \subset \mathbb{R}^N$ be an open set with uniformly Lipschitz boundary. Then there exists a continuous linear operator $\mathfrak{L}: W^{1,\infty}(\Omega) \to W^{1,\infty}(\mathbb{R}^N)$ such that $\mathfrak{L}(u)(\mathbf{x}) = u(\mathbf{x})$ for all $u \in W^{1,\infty}(\Omega)$.

The interconnection between stability of the nonlinear scheme (8) and stability of its scheme in variations is established by the following theorem. It is assumed that Ω_F in (8) need not be open. It is also assumed that the set Ω_F is bounded, i.e. there exists $r_0 < \infty$ such that $\Omega_F \subset B_0 \equiv B(0, r_0)$. Let $\mathbf{F} \equiv \{F_1, F_2, \ldots, F_N\}^T$ and let ∇F_i , $i = 1, 2, \ldots, N$, denote the distributional gradient of F_i . Definition of the distributional gradient (or distributional partial derivative) of F_i at $\mathbf{x} \in int(\Omega_F)$ can be found, e.g., in [31]. If $\mathbf{x} \in \partial \Omega_F$, while Ω_F is Lipschitz, then the distributional gradient, ∇F_i , is assumed to be equal to $\nabla \mathfrak{L}(F_i)(\mathbf{x})$.

Theorem 8 Consider the scheme (8). Let \mathbf{F} be bounded and let Ω_F be internal path-connected, bounded, and Lipschitz. Then, the scheme, (8), will be stable in terms of Definition 2 iff its scheme in variations, (13), will be stable.

Proof. Let the scheme (8) be stable, then, in view of Theorem 4, \mathbf{F} will be locally Lipschitz for a common constant C, and, hence, in view of Theorem 5, $F_i \in W^{1,\infty}(int(\Omega_F))$, $i=1,2,\ldots,N$. By virtue of Theorem 7 we can find $F_i^* \in W^{1,\infty}(\mathbb{R}^N)$ such that $F_i^*(\mathbf{x}) = F_i(\mathbf{x})$ on $int(\Omega_F)$. Hence $F_i^* \in W^{1,\infty}(B_0)$, where B_0 is an open ball such that $\Omega_F \subset B_0 \subset \mathbb{R}^N$. Notice, $\|\nabla F_i^*\|_{\infty}$, $i=1,2,\ldots,N$, is bounded on the ball B_0 and, naturally, on the domain $\Omega_F \subset B_0$,

since $F_i^* \in W^{1,\infty}(B_0)$. In such a case there exists a common constant C such that $\|\mathbf{F}'\| \leq C$ for all $\mathbf{x} \in \Omega_F$, where $\|\mathbf{F}'\| \equiv \|\mathbf{f}_F\|$, $\mathbf{f}_F \equiv \{\|\nabla F_1\|_{\infty}, \|\nabla F_2\|_{\infty}, \dots, \|\nabla F_N\|_{\infty}\}^T$. Hence, in view of Lemma 6, the scheme in variations, (13), is stable.

Conversely, suppose that the scheme in variations, (13), is stable. In view of Lemma 6, there exists a common constant C such that $\|\mathbf{F}'\| \leq C$ on Ω_F , and, hence, $\|\nabla F_i\|_{\infty}$, $i=1,2,\ldots,N$, is bounded on Ω_F . Consequently, $F_i \in W^{1,\infty}\left(int\left(\Omega_F\right)\right)$, $i=1,2,\ldots,N$. Since $int\left(\Omega_F\right)$ is an extension domain for $W^{1,\infty}\left(int\left(\Omega_F\right)\right)$, we can find, in view of Theorem 7, $F_i^* \in W^{1,\infty}\left(\mathbb{R}^N\right)$ such that $F_i^*(\mathbf{x}) = F_i(\mathbf{x})$ on $int\left(\Omega_F\right)$. In view of Theorem 5, F_i^* , $i=1,2,\ldots,N$, will be locally Lipschitz on \mathbb{R}^N (and, hence, on Ω_F) for a common constant. Then F_i , $i=1,2,\ldots,N$, and hence \mathbf{F} will be locally Lipschitz on Ω_F , since Ω_F being Lipschitz satisfies the cone property. Thus, in view of Theorem 4, the scheme (8) will be stable.

Notice, if **F** in the scheme (8) is Gateaux-differentiable, then ∇F_i (see Theorem 8) denotes the classical gradient of F_i , and, hence, it may be taken that $\mathbf{f}_F = \{\|\nabla F_1\|, \|\nabla F_2\|, \dots, \|\nabla F_N\|\}^T$.

Remark 9 Definition 1 is "restored in its rights" for domains being internal path-connected, bounded, and Lipschitz (see the proof of Theorem 8). Namely, considering the scheme (8) with such a domain, we find that the scheme will be stable in terms of Definition 2 iff it will be stable in terms of Definition 1.

3 Construction of stable central schemes

In this section we consider explicit schemes on a uniform grid with time step Δt and spatial mesh size Δx , as applied to 1-D hyperbolic equations. In view of the CFL condition [32], we assume for the explicit schemes that $\Delta t = O(\Delta x)$. Moreover, we will also assume that $\Delta x = O(\Delta t)$, since a central scheme generates, in general, a conditional approximation (see [8]). In such a case, the following inequalities will be valid, for sufficiently small Δt and Δx ,

$$\nu_0 \Delta t \le \Delta x \le \mu_0 \Delta t, \quad \nu_0, \mu_0 = const, \ 0 < \nu_0 \le \mu_0. \tag{21}$$

We will focus on the following 1-D version of the problem (1):

$$\frac{\partial \mathbf{u}}{\partial t} + \frac{\partial}{\partial x} \mathbf{f}(\mathbf{u}) = 0, \ t_n < t \le t_{n+1} \equiv t_n + \Delta t, \quad \mathbf{u}(x, t_n) = \mathbf{u}^n(x),$$
 (22)

Since the system (22) is hyperbolic, the Jacobian matrix of $\mathbf{f}(\mathbf{u})$ possesses M linearly independent eigenvectors (see, e.g., [18]). In addition, it is also assumed that

$$\sup_{\mathbf{u} \in \Omega_{\mathbf{u}}} \|\mathbf{A}\|_{2} \le \lambda_{\max} < \infty, \quad \mathbf{A} = \frac{\partial \mathbf{f}(\mathbf{u})}{\partial \mathbf{u}}.$$
 (23)

Using the central differencing, we write

$$\frac{\partial \mathbf{u}}{\partial t}\Big|_{t=t_{n+0.25}, \ x=x_{i+0.5}} = \frac{\mathbf{u}_{i+0.5}^{n+0.5} - \mathbf{u}_{i+0.5}^{n}}{0.5\Delta t} + O\left((\Delta t)^{2}\right), \tag{24}$$

$$\left. \frac{\partial \mathbf{f}}{\partial x} \right|_{t=t_{n+0.25}, \ x=x_{i+0.5}} = \frac{\mathbf{f}_{i+1}^{n+0.25} - \mathbf{f}_{i}^{n+0.25}}{\Delta x} + O\left((\Delta x)^{2}\right).$$
 (25)

By virtue of (24)-(25) we approximate (22) on the cell $[x_i, x_{i+1}] \times [t_n, t_{n+0.5}]$ by the following difference equation

$$\mathbf{v}_{i+0.5}^{n+0.5} = \mathbf{v}_{i+0.5}^{n} - \frac{\Delta t}{2\Delta x} \left(\mathbf{g}_{i+1}^{n+0.25} - \mathbf{g}_{i}^{n+0.25} \right). \tag{26}$$

As usual, the mathematical treatment for the second step (i.e., on the cell $[x_{i-0.5}, x_{i+0.5}] \times [t_{n+0.5}, t_{n+1}]$) of a staggered scheme will, in general, not be included in the text, because it is quite similar to the treatment for the first step.

Considering that (26) approximates (22) with the accuracy $O((\Delta x)^2 + (\Delta t)^2)$, the next problem is to approximate $\mathbf{v}_{i+0.5}^n$ and $\mathbf{g}_i^{n+0.25}$ in such a way as to retain the accuracy of the approximation. For instance, the following approximations

$$\mathbf{v}_{i+0.5}^{n} = 0.5\left(\mathbf{v}_{i}^{n} + \mathbf{v}_{i+1}^{n}\right) + O\left(\left(\Delta x\right)^{2}\right), \quad \mathbf{g}_{i}^{n+0.25} = \mathbf{f}\left(\mathbf{v}_{i}^{n}\right) + O\left(\Delta t\right), \quad (27)$$

leads to the staggered form of the famed LxF scheme that is of the first-order approximation (see, e.g., [18, p. 170]). One way to obtain a higher-order scheme is to use a higher order interpolation. At the same time it is required of the interpolant to be monotonicity preserving. Notice, the classic cubic spline [46] does not possess such a property (see Figure 1a). In the following, subsection 3.1, we consider the problem of high-order interpolation of $\mathbf{v}_{i+0.5}^n$ in (26) with closer inspection.

3.1 Monotone C^1 piecewise cubics in construction of explicit schemes

We will consider some theoretical aspects for high-order interpolation and employment of monotone C^1 piecewise cubics (see, e.g., [14], [27]) in construction of monotone explicit schemes.

Let $\mathbf{p} = \mathbf{p}(x) \equiv \{p^1(x), \dots, p^k(x), \dots, p^m(x)\}^T$ be an interpolant, and let

$$\mathbf{p}_{i} = \mathbf{p}\left(x_{i}\right), \quad \mathbf{p}_{i}' = \mathbf{p}'\left(x_{i}\right), \quad \Delta\mathbf{p}_{i} = \mathbf{p}_{i+1} - \mathbf{p}_{i},$$

$$\mathbf{p}_{i}' = \mathbb{A}_{i} \cdot \frac{\Delta \mathbf{p}_{i}}{\Delta x}, \quad \mathbf{p}_{i+1}' = \mathbb{B}_{i} \cdot \frac{\Delta \mathbf{p}_{i}}{\Delta x},$$
 (28)

where \mathbf{p}'_i denotes the derivative of the interpolant at $x = x_i$. The diagonal matrices \mathbb{A}_i and \mathbb{B}_i in (28) are defined as follows

$$\mathbb{A}_{i} = diag\left\{\alpha_{i}^{1}, \alpha_{i}^{2}, \dots, \alpha_{i}^{m}\right\}, \ \mathbb{B}_{i} = diag\left\{\beta_{i}^{1}, \beta_{i}^{2}, \dots, \beta_{i}^{m}\right\}.$$
 (29)

If $\mathbf{p} = \mathbf{p}(x)$ is a C^1 piecewise cubic interpolant (see, e.g., [14], [27]), then it will be component-wise monotone on $[x_i, x_{i+1}]$ iff one of the following conditions (see [14], [27]) is satisfied:

$$(\alpha_i^k - 1)^2 + (\alpha_i^k - 1)(\beta_i^k - 1) + (\beta_i^k - 1)^2 - 3(\alpha_i^k + \beta_i^k - 2) \le 0, \quad (30)$$

$$\alpha_i^k + \beta_i^k \le 3, \quad \alpha_i^k \ge 0, \ \beta_i^k \ge 0, \quad \forall i, k.$$
 (31)

The region of monotonicity is shown in Figure 1b. The results of implementing a monotone C^1 piecewise cubic interpolation when compared with the classic cubic spline interpolation, are depicted in Figure 1a.

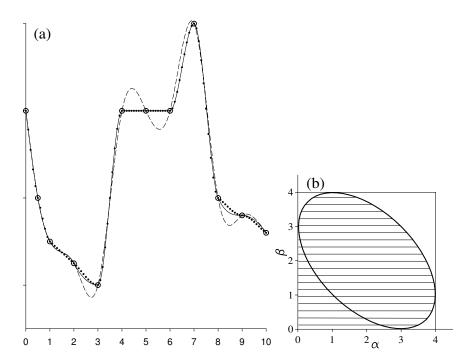


Figure 1: Monotone piecewise cubic interpolation. (a) Interpolation of a 1-D tabulated function. Circles: prescribed tabulated values; Dashed line: classic cubic spline; Solid line: monotone piecewise cubic under $\alpha, \beta \leq 3$; Dotted line: monotone piecewise cubic under $\alpha, \beta \leq 1$. (b) Necessary and sufficient conditions for monotonicity. Horizontal hatching: region of monotonicity; Unshaded: cubic is non-monotone.

We note (Figure 1a) that the constructed function produces monotone interpolation and has a small, or even practically zero, deviation from the classic cubic spline at some sections where the classic cubic spline is monotone.

Let us consider the problem of monotone high-order approximation of $\mathbf{v}_{i+0.5}^n$ in (26). First, given the values, \mathbf{p}_i and \mathbf{p}_{i+1} , of the interpolant $\mathbf{p} = \mathbf{p}(x)$ and the approximate estimates, \mathbf{d}_i and \mathbf{d}_{i+1} , of its derivatives, namely

$$\mathbf{p}_{i}' = \mathbf{d}_{i} + O\left(\left(\Delta x\right)^{s}\right),\tag{32}$$

we obtain the following interpolation formula

$$\mathbf{p}_{i+0.5} = 0.5 \left(\mathbf{p}_i + \mathbf{p}_{i+1} \right) - \frac{\Delta x}{8} \left(\mathbf{d}_{i+1} - \mathbf{d}_i \right) + O\left(\left(\Delta x \right)^r \right), \tag{33}$$

where

$$r = \min\left(4, s + 1\right). \tag{34}$$

Actually, let $\mathbf{p}(x)$ be sufficiently smooth, then, using the Taylor series expansion of the function $\mathbf{p}(x)$, we obtain:

$$\mathbf{p}_{i+1} + \mathbf{p}_i = 2\mathbf{p}_{i+0.5} + \mathbf{p}_{i+0.5}'' \left(\frac{\Delta x}{2}\right)^2 + \frac{\mathbf{p}_{i+0.5}''''}{12} \left(\frac{\Delta x}{2}\right)^4 + O\left((\Delta x)^6\right).$$
 (35)

In a similar manner, using the Taylor series expansion of the function $\mathbf{p}'(x)$, we can show that

$$\mathbf{p}'_{i+1} - \mathbf{p}'_{i} = \mathbf{p}''_{i+05} \Delta x + \frac{\mathbf{p}'''_{i+05}}{3} \left(\frac{\Delta x}{2}\right)^{3} + O\left((\Delta x)^{5}\right)$$
(36)

By virtue of (36) we obtain from (35) that

$$\mathbf{p}_{i+0.5} = 0.5 \left(\mathbf{p}_i + \mathbf{p}_{i+1} \right) - \frac{\Delta x}{8} \left(\mathbf{p}'_{i+1} - \mathbf{p}'_i \right) + \frac{\mathbf{p}''''_{i+0.5}}{24} \left(\frac{\Delta x}{2} \right)^4 + O\left((\Delta x)^6 \right)$$
(37)

In view of (32) we obtain from (37) the following interpolation formula

$$\mathbf{p}_{i+0.5} = 0.5 \left(\mathbf{p}_i + \mathbf{p}_{i+1} \right) - \frac{\Delta x}{8} \left(\mathbf{d}_{i+1} - \mathbf{d}_i \right) + O\left((\Delta x)^4 + (\Delta x)^{s+1} \right).$$
 (38)

Hence, by virtue of (38) and (33), we conclude that (34) is, in general, valid. It is clear that if (32) is valid and $\mathbf{d}_{i+1} - \mathbf{d}_i = (\mathbf{p}'_{i+1} - \mathbf{p}'_i) + O((\Delta x)^{s+1})$, then $r = \min(4, s+2)$ in (33). Thus, if $\mathbf{p}(x)$ has a continuous fourth derivative and \mathbf{p}'_i can be estimated with accuracy $O((\Delta x)^3)$ (see, e.g., [26, p. 112]), then r = 4 in (33).

Let us note that instead of point values employed in the construction of the scheme (26), it can be used the cell averages (see, e.g., [5], [29], [32]). Then, given the values, $\overline{\mathbf{p}}_i$, of the averaged interpolant $\overline{\mathbf{p}} = \overline{\mathbf{p}}(x)$ and the estimates, $\overline{\mathbf{d}}_i$ and $\overline{\mathbf{d}}_{i+1}$, of its derivatives, namely $\overline{\mathbf{p}}_i' = \overline{\mathbf{d}}_i + O\left((\Delta x)^s\right)$ we obtain the same formula (33) written for the averaged values. Further, to estimate the flux at x_i , the point value \mathbf{p}_i can be reconstructed using a piecewise polynomial interpolation (e.g., [5], [29]). Notice, we need not any reconstruction when dealing with first- and second-order schemes, since the point values agree with the corresponding cell averages to $O((\Delta x)^2)$ (e.g., [29], [32]). We will therefore omit the bar notation when dealing with such schemes.

Instead of (33), we will employ the following interpolation formula

$$\mathbf{p}_{i+0.5} = 0.5 \left(\mathbf{p}_i + \mathbf{p}_{i+1} \right) - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1} - \mathbf{d}_i \right), \quad 0 \le \varkappa \le 1,$$
 (39)

where, for some reason, a new parameter \varkappa is used. Consider, for instance, the case when the point value $\mathbf{v}_{i+0.5}^n$ in the scheme (26) is estimated by the cell average calculated on the basis of the monotone C^1 piecewise cubics. Note that this estimation is of second-order. Such an approach leads to the interpolation of the point value by formula (39) with the parameter $\varkappa = 2/3$. Obviously, the interpolation (39) is of second-order under $0 \le \varkappa < 1$, and (39) coincides with (33) under $\varkappa = 1$ to give fourth order accuracy. Notice, if $0 < \varkappa < 1$, then $\mathbf{p}_{i+0.5}$ can be estimated with accuracy $O((\Delta x)^4)$ by (39) provided that $1 - \varkappa = O((\Delta x)^2)$. Actually, by virtue of (36) and (32), we obtain that

$$\varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1} - \mathbf{d}_i \right) = \frac{\Delta x}{8} \left(\mathbf{d}_{i+1} - \mathbf{d}_i \right) + (1 - \varkappa) \frac{\left(\Delta x \right)^2}{8} \frac{\mathbf{d}_{i+1} - \mathbf{d}_i}{\Delta x} = \frac{\Delta x}{8} \left(\mathbf{d}_{i+1} - \mathbf{d}_i \right)$$

$$+ \frac{O((\Delta x)^4)}{8} \mathbf{p}_{i+05}'' + O\left((\Delta x)^{s+3}\right) = \frac{\Delta x}{8} \left(\mathbf{d}_{i+1} - \mathbf{d}_i\right) + O((\Delta x)^4). \tag{40}$$

The approximation of derivatives \mathbf{p}'_i can be done by the following three steps [14]: (i) initialization of the derivatives $\mathbf{p}'_i = \mathbf{d}_i + O((\Delta x)^s)$; (ii) choice of subregion of monotonicity; (iii) modification of the initialized derivatives, \mathbf{d}_i , to produce a monotone interpolant.

The matter of initialization of the derivatives is the most subtle issue of this algorithm. Thus, using the two-point or the three-point (centered) difference formula (e.g. [27], [41]) we obtain, in general, second-order approximation. Performing the initialization of the derivatives $\mathbf{p}'_i = \mathbf{d}_i + O\left((\Delta x)^s\right)$ in the interpolation formula (39) by the classic cubic spline [46] interpolation, we obtain s = 3 (e.g., [26], [27]). The same accuracy can be achieved by using the four-point approximation [27]. However, the efficiency of the algorithm based on the classic cubic spline interpolation is comparable with the one based on the four-point approximation, as the number of multiplications and divisions per one node is approximately the same for both algorithms.

Obviously, for each interval $[x_i, x_{i+1}]$ in which the initialized derivatives \mathbf{d}_i , \mathbf{d}_{i+1} such that at least one point (α_i^k, β_i^k) does not belong to the region of monotonicity (30)-(31), the derivatives \mathbf{d}_i , \mathbf{d}_{i+1} must be modified to $\widetilde{\mathbf{d}}_i$, $\widetilde{\mathbf{d}}_{i+1}$ such that the point $(\widetilde{\alpha}_i^k, \widetilde{\beta}_i^k)$ will be in the region of monotonicity. Moreover, it is desirable to modify the derivatives $\mathbf{d} \equiv \left\{ \dots, \mathbf{d}_{i-1}^T, \mathbf{d}_i^T, \mathbf{d}_{i+1}^T, \dots \right\}^T$ in such a way that $\widetilde{\mathbf{d}} \equiv \left\{ \dots, \widetilde{\mathbf{d}}_{i-1}^T, \widetilde{\mathbf{d}}_i^T, \widetilde{\mathbf{d}}_{i+1}^T, \dots \right\}^T$ would be the solution to the following mathematical programming problem:

$$\left\| \mathbf{d} - \widetilde{\mathbf{d}} \right\| \to \min_{\widetilde{\mathbf{d}}}.$$
 (41)

The modification of the initialized derivatives, would be much simplified if we take a square as a subregion of monotonicity. In connection with this, we will make use the subregions of monotonicity represented in the following form:

$$0 \le \alpha_i^k \le 4\aleph, \quad 0 \le \beta_i^k \le 4\aleph, \quad \forall i, k, \tag{42}$$

where \aleph is a monotonicity parameter. If $\aleph = 0.75$, then the subregion (42), called de Boor-Swartz box [27], coincides with the one used by Fritsch and Carlson [14]. It is proven in [7] that the interpolation (39) will be monotone, i.e. the value of an arbitrary component of $\mathbf{p}_{i+0.5}$ will be between the corresponding components of \mathbf{p}_i and \mathbf{p}_{i+1} , iff (42) will be valid provided that $0 \le \aleph \le 1$. Notice, in such a case the point (α_i^k, β_i^k) may be out of the region of monotonicity shown in Figure 1b.

To fulfill the conditions of monotonicity (42), the modification of derivatives $\mathbf{d}_i = \{d_i^1, d_i^2, \dots, d_i^m\}$ can be done by the following algorithm suggested, in fact, by Fritsch and Carlson [14] (see also [27]):

$$S_i^k := 4 \aleph \min \operatorname{mod}(\Delta_{i-1}^k, \Delta_i^k), \quad \widetilde{d}_i^k := \min \operatorname{mod}(d_i^k, S_i^k), \quad \aleph = const, \quad (43)$$

where $\Delta_i^k = (p_{i+1}^k - p_i^k) / \Delta x$, the function min mod(x, y) is defined (e.g., [27], [29], [36], [41], [50]) as follows

$$\min \operatorname{mod}(x, y) \equiv \frac{1}{2} \left[sgn(x) + sgn(y) \right] \min \left(|x|, |y| \right). \tag{44}$$

3.2 Construction of nonstaggered central schemes

We will consider explicit central schemes on a uniform grid with time step Δt and spatial mesh size Δx . In view of the interpolation formula (39), the staggered scheme (26) becomes

$$\mathbf{v}_{i+0.5}^{n+0.5} = 0.5 \left(\mathbf{v}_{i+1}^{n} + \mathbf{v}_{i}^{n} \right) - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n} - \mathbf{d}_{i}^{n} \right) - \frac{\Delta t}{2} \frac{\mathbf{f} \left(\mathbf{v}_{i+1}^{n} \right) - \mathbf{f} \left(\mathbf{v}_{i}^{n} \right)}{\Delta x}, \quad (45)$$

where \mathbf{d}_i^n denotes the derivative of the interpolant at $x = x_i$, the range of values for the parameter \varkappa is the segment $0 \le \varkappa \le 1$. If $\varkappa = 0$, then the scheme (45) coincides with the LxF scheme. It is shown in [7] that the first-order, $O(\Delta t + (\Delta x)^r / \Delta t + (\Delta x)^2)$, scheme (45) generates a conditional approximation, because it approximates (22) only if $(\Delta x)^r / \Delta t \to 0$ as $\Delta x \to 0$ and $\Delta t \to 0$, where r is the order of approximation of $\mathbf{v}_{i+0.5}^n$ by the interpolation formula (39). The scheme (45) is abbreviated in [8] as COS1.

Let us develop new nonoscillatory central schemes, which are based on regular, nonstaggered spatial grids. Using the central differencing, we approximate (22) on the cell $[x_{i-1}, x_{i+1}] \times [t_n, t_{n+1}]$ by the following difference equation

$$\mathbf{v}_{i}^{n+1} = \mathbf{v}_{i}^{n} - \frac{\Delta t}{2\Delta x} \left(\mathbf{g}_{i+1}^{n+0.5} - \mathbf{g}_{i-1}^{n+0.5} \right). \tag{46}$$

Notice, the approximation $\mathbf{g}_{i}^{n+0.5} = \mathbf{f}(\mathbf{v}_{i}^{n}) + O(\Delta t)$, leads to the first-order, $O(\Delta t + (\Delta x)^{2})$, scheme

$$\mathbf{v}_{i}^{n+1} = \mathbf{v}_{i}^{n} - \frac{\Delta t}{2\Delta x} \left(\mathbf{f}_{i+1}^{n} - \mathbf{f}_{i-1}^{n} \right), \tag{47}$$

which is absolutely unstable [16, p. 113], [32, p. 71]. Let us construct a stable scheme on the basis of (46). The interpolation formula (39), as applied to the cell $[x_{i-1}, x_{i+1}] \times [t_n, t_{n+1}]$, becomes

$$\mathbf{p}_{i} = 0.5 \left(\mathbf{p}_{i-1} + \mathbf{p}_{i+1} \right) - \varkappa \frac{\Delta x}{4} \left(\mathbf{d}_{i+1} - \mathbf{d}_{i-1} \right).$$
 (48)

Equality (48), after elementary transformations, leads to

$$\mathbf{p}_{i} = 0.25 \left(\mathbf{p}_{i-1} + 2\mathbf{p}_{i} + \mathbf{p}_{i+1} \right) - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1} - \mathbf{d}_{i-1} \right). \tag{49}$$

Using (49) we obtain from (47) the following nonstaggered central scheme

$$\mathbf{v}_{i}^{n+1} = \frac{\mathbf{v}_{i+1}^{n} + 2\mathbf{v}_{i}^{n} + \mathbf{v}_{i-1}^{n}}{4} - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n} - \mathbf{d}_{i-1}^{n} \right) - \frac{\Delta t}{2\Delta x} \left(\mathbf{f}_{i+1}^{n} - \mathbf{f}_{i-1}^{n} \right). \tag{50}$$

where $0 \le \varkappa \le 1$. In view of (32)-(34), as applied to the cell $[x_{i-1}, x_{i+1}] \times [t_n, t_{n+1}]$, the local truncation error, ψ , on a sufficiently smooth solution $\mathbf{u}(x, t)$ to (22) is found to be

$$\psi = O(\Delta t + (\Delta x)^2) + O\left(\frac{(\Delta x)^r}{\Delta t}\right) + (1 - \varkappa)O\left(\frac{(\Delta x)^2}{\Delta t}\right),\tag{51}$$

where $r = \min(4, s+1)$, s denotes the order of approximation in (32). Note that the scheme (50) will be $O(\Delta t + (\Delta x)^2)$ accurate if $r \geq 3$ as well as $1 - \varkappa = O(\Delta x)$. The first-order nonstaggered central scheme (50), being a monotone approximation of the 1-D equation (22) by C^1 cubics, will be abbreviated to MAC1.

Let us consider the transformation from the scheme (47) to (50) with closer inspection. We will use the, so called, first differential approximation of the scheme (47) ([15, p. 45], [49, p. 376]; see also 'modified equations' in [15, p. 45], [32], [35]). As reported in [15], [49], this heuristic method was originally presented by Hirt (1968) (see [15, p. 45]) as well as by Shokin and Yanenko (1968) (see [49, p. 376]), and has since been widely employed in the development of stable difference schemes for PDEs.

The first differential approximation of the scheme (47) is the following:

$$\frac{\partial \mathbf{u}}{\partial t} + \frac{\partial}{\partial x} \mathbf{f} \left(\mathbf{u} \right) = -\frac{\Delta t}{2} \frac{\partial}{\partial x} \left(\mathbf{A}^2 \cdot \frac{\partial \mathbf{u}}{\partial x} \right), \quad \mathbf{A}^2 \equiv \frac{\partial \mathbf{f}}{\partial \mathbf{u}} \cdot \frac{\partial \mathbf{f}}{\partial \mathbf{u}}. \tag{52}$$

The negative diffusion term in (52) is associated with instability, as this therm is the source of energy resulting in an unlimited growth of the amplitude of the solution (see, e.g., [51], [52]). Notice, for the sake of convenience we use the same notation in (22) and in (52) in spite of the fact that these equations are different. We will use such an approach if it does not lead to confusion.

To obtain (50), we, in fact, added to the right-hand side of (47) the value

$$\mathbf{E}_{i} \equiv \frac{\mathbf{v}_{i+1}^{n} - 2\mathbf{v}_{i}^{n} + \mathbf{v}_{i-1}^{n}}{4} - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n} - \mathbf{d}_{i-1}^{n} \right). \tag{53}$$

Let the interpolant $\mathbf{p} = \mathbf{p}(x)$ be such that $\mathbf{p}_i = \mathbf{v}_i^n$. If $\varkappa = 1$ and $\mathbf{d}_i = \mathbf{p}_i'$, then, by virtue of (32)-(37) as applied to the cell $[x_{i-1}, x_{i+1}] \times [t_n, t_{n+1}]$, we find that

$$\mathbf{E}_{i} = -\frac{\left(\Delta x\right)^{4}}{48}\mathbf{p}_{i}^{\prime\prime\prime\prime} + O\left(\left(\Delta x\right)^{6}\right). \tag{54}$$

Thus, when we add \mathbf{E}_i of (53) to the right-hand side of (47), we, in fact, add the negative fourth-order diffusion term of (54) to the right-hand side of the first differential approximation (52). Such a term, in contrast to a negative second-order diffusion term, stabilizes the system and produces dissipative effect (see, e.g., [51], [52]).

If $\varkappa = 1$ and \mathbf{p}'_i is estimated by the central-difference derivative, i.e. $\mathbf{d}_i = (\mathbf{p}_{i+1} - \mathbf{p}_{i-1}) / (2\Delta x)$, then

$$\mathbf{p}_{i}' = \mathbf{d}_{i} - \frac{\left(\Delta x\right)^{2}}{6} \mathbf{p}_{i}^{"'} + O\left(\left(\Delta x\right)^{4}\right). \tag{55}$$

In such a case we obtain, instead of (54), that

$$\mathbf{E}_{i} = -\frac{3\left(\Delta x\right)^{4}}{48}\mathbf{p}_{i}^{\prime\prime\prime\prime} + O\left(\left(\Delta x\right)^{5}\right). \tag{56}$$

We can see from (56) that the estimation of the derivative \mathbf{p}'_i by the central-difference derivative leads to the same order of accuracy as it was in the previous example when $\mathbf{d}_i = \mathbf{p}'_i$. Furthermore, the scheme (50) will be more dissipative in such a case, since the coefficient of the fourth-order diffusion term in (56) is three times bigger than the similar coefficient in (54).

If $0 \le \varkappa \le 1$ and \mathbf{p}_i' is estimated by the central-difference derivative, then we find

$$\mathbf{E}_{i} = (1 - \varkappa) \frac{(\Delta x)^{2}}{4} \mathbf{p}_{i}^{"} + (1 - 4\varkappa) \frac{(\Delta x)^{4}}{48} \mathbf{p}_{i}^{""} + O\left((\Delta x)^{5}\right).$$
 (57)

Note that the estimations (54), (56), and (57) will be valid if the derivative \mathbf{d}_i need not be modified, i.e. (42) will be valid without implementing the algorithm (43). Otherwise, there is a risk to add negative second-order diffusion terms to the right-hand side of (52). Let us consider, for instance, the case when $\varkappa = 1$ and the derivative \mathbf{p}'_{i+1} is estimated by the right-difference derivative, i.e. $\mathbf{d}_{i+1} = (\mathbf{p}_{i+2} - \mathbf{p}_{i+1})/\Delta x$, however \mathbf{p}'_{i-1} is estimated by the central-difference derivative. In such a case we obtain:

$$\mathbf{E}_{i} = -\frac{(\Delta x)^{2}}{16}\mathbf{p}_{i+1}^{"} - \frac{1.25(\Delta x)^{4}}{16}\mathbf{p}_{i+1}^{""} + O\left((\Delta x)^{5}\right).$$
 (58)

In view of (58) and (57), we conclude that the part played by the parameter \varkappa in suppressing spurious oscillations produced by negative second-order diffusion terms can be very important.

Interestingly, there is a possibility to improve the scheme (45) by introducing an additional positive numerical viscosity, i.e., using the vanishing viscosity method [18], [32], such that the scheme's order of accuracy would increase up to $O((\Delta t)^2 + (\Delta x)^2)$. Such an approach [8] leads to the following second order central scheme

$$\mathbf{v}_{i+0.5}^{n+0.5} = 0.5 \left(\mathbf{v}_{i+1}^{n} + \mathbf{v}_{i}^{n} \right) - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n} - \mathbf{d}_{i}^{n} \right) - \frac{\Delta t}{2} \frac{\mathbf{f} \left(\mathbf{v}_{i+1}^{n} \right) - \mathbf{f} \left(\mathbf{v}_{i}^{n} \right)}{\Delta x} + \frac{\left(\Delta t \right)^{2}}{8 \Delta x} \left[\left(\mathbf{A}_{i+1}^{n} \right)^{2} \cdot \mathbf{d}_{i+1}^{n} - \left(\mathbf{A}_{i}^{n} \right)^{2} \cdot \mathbf{d}_{i}^{n} \right], \quad \mathbf{A} \equiv \frac{\partial \mathbf{f}}{\partial \mathbf{u}},$$
 (59)

where \mathbf{d}_{i}^{n} is the derivative of the interpolant at $x = x_{i}$. The scheme (59) is abbreviated in [8] as COS2.

Let us develop a nonstaggered central scheme that will be of second order. Using the central differencing, we approximate (22) at the point $x = x_i$, $t = t_{n+05}$ by the following equation:

$$\mathbf{v}_{i}^{n+1} = \mathbf{v}_{i}^{n} - \Delta t \left. \frac{\partial \mathbf{f}}{\partial x} \right|_{x=x_{i}, \ t=t_{n+05}}.$$
 (60)

Using Taylor series expansion, we approximate $\mathbf{g}_i^{n+0.5}$ in (46) with the accuracy $O((\Delta t)^2)$:

$$\mathbf{g}_{i}^{n+0.5} = \mathbf{f}\left(\mathbf{v}_{i}^{n}\right) + \left. \frac{\partial \mathbf{f}\left(\mathbf{v}_{i}^{n}\right)}{\partial t} \right|_{t=t_{n}} \frac{\Delta t}{2} + O\left(\Delta t^{2}\right). \tag{61}$$

By virtue of the PDE system, (22), we find

$$\frac{\partial \mathbf{f}}{\partial t} = \frac{\partial \mathbf{f}}{\partial \mathbf{u}} \cdot \frac{\partial \mathbf{u}}{\partial t} = -\frac{\partial \mathbf{f}}{\partial \mathbf{u}} \cdot \frac{\partial \mathbf{f}}{\partial \mathbf{u}} \cdot \frac{\partial \mathbf{u}}{\partial x} = -\left(\frac{\partial \mathbf{f}}{\partial \mathbf{u}}\right)^2 \cdot \frac{\partial \mathbf{u}}{\partial x}.$$
 (62)

By virtue of (49) and (61)-(62), we find

$$\mathbf{v}_{i}^{n+1} = \frac{\mathbf{v}_{i-1}^{n} + 2\mathbf{v}_{i}^{n} + \mathbf{v}_{i+1}^{n}}{4} - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n} - \mathbf{d}_{i-1}^{n} \right) - \Delta t \left. \frac{\partial \mathbf{f}}{\partial x} \right|_{x=x_{i}, t=t_{n}} + \frac{(\Delta t)^{2}}{2} \left. \frac{\partial}{\partial x} \left(\mathbf{A}^{2} \frac{\partial \mathbf{u}}{\partial x} \right) \right|_{x=x_{i}, t=t_{n}}, \quad \mathbf{A}^{2} \equiv \frac{\partial \mathbf{f}}{\partial \mathbf{u}} \cdot \frac{\partial \mathbf{f}}{\partial \mathbf{u}}.$$
(63)

Then, approximating the last two terms in the right-hand side of (63) with the second order accuracy, we obtain

$$\mathbf{v}_{i}^{n+1} = \frac{\mathbf{v}_{i-1}^{n} + 2\mathbf{v}_{i}^{n} + \mathbf{v}_{i+1}^{n}}{4} - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n} - \mathbf{d}_{i-1}^{n} \right) - \frac{\Delta t}{2\Delta x} \left(\mathbf{f}_{i+1}^{n} - \mathbf{f}_{i-1}^{n} \right)$$

$$+\frac{\varsigma\left(\Delta t\right)^{2}}{2\left(\Delta x\right)^{2}}\left[\mathbf{D}_{i+0.5}^{n}\cdot\left(\mathbf{v}_{i+1}^{n}-\mathbf{v}_{i}^{n}\right)-\mathbf{D}_{i-0.5}^{n}\cdot\left(\mathbf{v}_{i}^{n}-\mathbf{v}_{i-1}^{n}\right)\right],\ \varsigma=const\geq0,\ (64)$$

where $\mathbf{D}_{i+0.5}^n = 0.5 \left[(\mathbf{A}_i^n)^2 + (\mathbf{A}_{i+1}^n)^2 \right]$, the parameter ς is introduced by analogy with \varkappa in (45). Scheme (64) coincides with the scheme MAC1, (50), provided that $\varsigma = 0$. The last term in right-hand side of (64) can be seen as the

non-negative numerical viscosity introduced into the first order scheme (50). Thus, we are dealing with the vanishing viscosity method [18], [32] and, hence, in view of [18, Theorem 3.3], the scheme, (64), satisfies the entropy condition. To make this point a little bit more clear, we consider a viscous perturbation of the system (22). We associate with (22) the following parabolic system

$$\frac{\partial \mathbf{u}_{\epsilon}}{\partial t} + \frac{\partial}{\partial x} \mathbf{f} \left(\mathbf{u}_{\epsilon} \right) = \epsilon \frac{\partial}{\partial x} \left(\mathbf{A}^{2} \cdot \frac{\partial \mathbf{u}_{\epsilon}}{\partial x} \right), \quad \epsilon > 0, \tag{65}$$

where the right-hand side can be viewed as a viscosity term. It is assumed that Mock's assumption of admissibility [18, p. 32] is valid, i.e. the matrix $U'' \cdot \mathbf{A}^2$ is positive-definite. Here $U = U\left(\mathbf{u}_{\epsilon}\right)$ denotes a strictly convex entropy, U'' denotes the Hessian matrix of U. We recover solutions of (22) as the limits of solutions to (65) as $\epsilon \to 0$. If we divide the scheme COS1, (50), by Δt and group all the terms of this scheme together in its left-hand side, then this group, in view of the above, is the finite difference approximation of the left-hand side of Eq. (65) on the cell $[x_{i-1}, x_{i+1}] \times [t_n, t_{n+1}]$. The right-hand side of Eq. (65) is approximated as following

$$\epsilon \frac{\partial}{\partial x} \left(\mathbf{A}^2 \cdot \frac{\partial \mathbf{u}_{\epsilon}}{\partial x} \right) \approx \epsilon \frac{\mathbf{D}_{i+0.5}^n \cdot \left(\mathbf{v}_{i+1}^n - \mathbf{v}_i^n \right) - \mathbf{D}_{i-0.5}^n \cdot \left(\mathbf{v}_i^n - \mathbf{v}_{i-1}^n \right)}{\left(\Delta x \right)^2}, \tag{66}$$

where $\mathbf{D}_{i\pm0.5}^n=0.5\left[\left(\mathbf{A}_i^n\right)^2+\left(\mathbf{A}_{i\pm1}^n\right)^2\right]$. If $\epsilon=\varsigma\Delta t/2$, then we obtain the scheme (64) as the approximation of Eq. (65) on the cell $[x_{i-1},x_{i+1}]\times[t_n,t_{n+1}]$. Thus, owing to the last term in the right-hand side of (64), the grid function \mathbf{v}_i^n in (64) can be considered as Lipschitz-continuous. Moreover, owing to this term, the scheme (64) is $O((\Delta x)^2+(\Delta t)^2)$ accurate provided that $\varkappa=\varsigma=1$. The nonstaggered central scheme (64), being a monotone approximation of the 1-D equation (22) with the second order, will be abbreviated to MAC2.

3.3 Stability of the developed schemes

It is proven in [8] that the scheme COS2, (59), will be, in general, stable if

$$\left|\varkappa - \varsigma C_r^2\right| \aleph + C_r \le 1, \quad C_r = \frac{\Delta t \lambda_{\text{max}}}{\Delta x}.$$
 (67)

Let us note that the stability condition (67) is not exactly correct, namely, there exist infrequent situations when (67) is valid nevertheless the scheme is not stable provided $\varkappa, \varsigma, \aleph = const.$ In particular, there could be a grid node (x_*, t_{**}) where the following inequality must be, but not valid:

$$\varsigma \aleph C_r^2 + C_r \le 1, \quad C_r = \frac{\Delta t \lambda_{\text{max}}}{\Delta r}.$$
 (68)

However, there is a possibility to take the parameters $\varkappa = \varkappa(x,t)$, $\varsigma = \varsigma(x,t)$, and $\aleph = \aleph(x,t)$ such that the scheme COS2, (59), will be stable under (67).

For instance, if we take $\varkappa, \varsigma, C_r = 1$, $\aleph = 0.5$ at $(x, t) \neq (x_*, t_{**})$ and $\aleph = 0$ at $(x, t) = (x_*, t_{**})$, then (68) will be also valid.

In this subsection we will use analogous technic [8] to prove the stability of the developed nonstaggered central schemes. In view of Theorem 8, the stability of the scheme MAC2, (64), will be investigated on the basis of its variational scheme. It is assumed that the bounded operator $\mathbf{A} = \partial \mathbf{f}(\mathbf{u}) / \partial \mathbf{u}$ in (22) is Fréchet-differentiable on the set $\Omega_{\mathbf{u}} \subset \mathbb{R}^M$, and its derivative is bounded on $\Omega_{\mathbf{u}}$. Considering that \mathbf{v}_i^n in (59) is Lipschitz-continuous, we write

$$\left\|\mathbf{v}_{i+1}^{n} - \mathbf{v}_{i}^{n}\right\|_{2} \le C_{v} \Delta x, \quad C_{v} = const. \tag{69}$$

By virtue of (28), the second term in right-hand side of (64) can be written in the form

$$\varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^n - \mathbf{d}_{i-1}^n \right) = \frac{\varkappa}{8} \left[\mathbb{B}_i^n \cdot \left(\mathbf{v}_{i+1}^n - \mathbf{v}_i^n \right) - \mathbb{A}_{i-1}^n \cdot \left(\mathbf{v}_i^n - \mathbf{v}_{i-1}^n \right) \right]. \tag{70}$$

Then, the variational scheme corresponding to (64) is the following

$$\delta \mathbf{v}_{i}^{n+1} - 0.25 \left(\delta \mathbf{v}_{i+1}^{n} + 2 \delta \mathbf{v}_{i}^{n} + \delta \mathbf{v}_{i-1}^{n} \right)$$

$$+ \frac{\Delta t}{2\Delta x} \left(\mathbf{A}_{i+1}^{n} \cdot \delta \mathbf{v}_{i+1}^{n} - \mathbf{A}_{i-1}^{n} \cdot \delta \mathbf{v}_{i-1}^{n} \right)$$

$$+ \frac{\varkappa}{8} \left[\mathbb{B}_{i}^{n} \cdot \left(\delta \mathbf{v}_{i+1}^{n} - \delta \mathbf{v}_{i}^{n} \right) - \mathbb{A}_{i-1}^{n} \cdot \left(\delta \mathbf{v}_{i}^{n} - \delta \mathbf{v}_{i-1}^{n} \right) \right]$$

$$-\varsigma \frac{(\Delta t)^{2}}{2\Delta x^{2}} \left[\mathbf{D}_{i+0.5}^{n} \cdot \left(\delta \mathbf{v}_{i+1}^{n} - \delta \mathbf{v}_{i}^{n} \right) - \mathbf{D}_{i-0.5}^{n} \cdot \left(\delta \mathbf{v}_{i}^{n} - \delta \mathbf{v}_{i-1}^{n} \right) \right] =$$

$$+\varsigma \frac{(\Delta t)^{2}}{2\Delta x^{2}} \left[\left(\delta \mathbf{D}_{i+0.5}^{n} \right) \cdot \left(\mathbf{v}_{i+1}^{n} - \mathbf{v}_{i}^{n} \right) - \left(\delta \mathbf{D}_{i-0.5}^{n} \right) \cdot \left(\mathbf{v}_{i}^{n} - \mathbf{v}_{i-1}^{n} \right) \right]$$

$$- \frac{\varkappa}{8} \left[\left(\delta \mathbb{B}_{i}^{n} \right) \cdot \left(\mathbf{v}_{i+1}^{n} - \mathbf{v}_{i}^{n} \right) - \left(\delta \mathbb{A}_{i-1}^{n} \right) \cdot \left(\mathbf{v}_{i}^{n} - \mathbf{v}_{i-1}^{n} \right) \right]. \tag{71}$$

By virtue of (42) and (23), we find

$$\left\|\delta \mathbb{A}_{i}^{n}\right\|_{2}, \left\|\delta \mathbb{B}_{i}^{n}\right\|_{2} \leq 4\aleph, \quad \left\|\left(\mathbf{A}_{i}^{n}\right)^{2}\right\|_{2}, \left\|\left(\mathbf{A}_{i+1}^{n}\right)^{2}\right\|_{2} \leq \lambda_{\max}^{2}.$$
 (72)

Thus, we may write that

$$\|\delta \mathbf{D}_{i+0.5}^{n}\|_{2}, \|\delta \left[\mathbf{D}_{i-0.5}^{n}\right]\|_{2} \le 2\lambda_{\max}^{2}.$$
 (73)

Then, by virtue of (21), (69), (72)-(73), and since $0 \le \varkappa, \varsigma, \aleph \le 1$, we find the following estimate for the right-hand side of (71):

$$\varsigma \frac{\left(\Delta t\right)^{2}}{2\Delta x^{2}} \left\| \left(\delta \mathbf{D}_{i+0.5}^{n}\right) \cdot \left(\mathbf{v}_{i+1}^{n} - \mathbf{v}_{i}^{n}\right) - \left(\delta \mathbf{D}_{i-0.5}^{n}\right) \cdot \left(\mathbf{v}_{i}^{n} - \mathbf{v}_{i-1}^{n}\right) \right\|_{2} \\
+ \left\| \frac{\varkappa}{8} \left[\left(\delta \mathbb{B}_{i}^{n}\right) \cdot \left(\mathbf{v}_{i+1}^{n} - \mathbf{v}_{i}^{n}\right) - \left(\delta \mathbb{A}_{i-1}^{n}\right) \cdot \left(\mathbf{v}_{i}^{n} - \mathbf{v}_{i-1}^{n}\right) \right] \right\|_{2} \le$$

$$\varsigma \frac{(\Delta t)^{2}}{2\Delta x^{2}} \left(\left\| \delta \mathbf{D}_{i+0.5}^{n} \right\|_{2} \left\| \mathbf{v}_{i+1}^{n} - \mathbf{v}_{i}^{n} \right\|_{2} + \left\| \delta \mathbf{D}_{i-0.5}^{n} \right\|_{2} \cdot \left\| \mathbf{v}_{i}^{n} - \mathbf{v}_{i-1}^{n} \right\|_{2} \right)
+ \frac{1}{8} \left(\left\| \delta \mathbb{B}_{i}^{n} \right\|_{2} \cdot \left\| \mathbf{v}_{i+1}^{n} - \mathbf{v}_{i}^{n} \right\|_{2} + \left\| \delta \mathbb{A}_{i-1}^{n} \right\|_{2} \cdot \left\| \mathbf{v}_{i}^{n} - \mathbf{v}_{i-1}^{n} \right\|_{2} \right) \le
\left(1 + 2C_{r}^{2} \right) C_{v} \mu_{0} \Delta t, \quad C_{r} = \frac{\Delta t \lambda_{\text{max}}}{\Delta x}.$$
(74)

Since the uniform stability with respect to the initial data implies the stability of scheme [48, pp. 390-392] (see also [8, Sec. 5]), we conclude, in view of (74), that the scheme (71) will be stable if the following scheme (i.e. (71) without the right-hand side) will be stable

$$\delta \mathbf{v}_{i}^{n+1} = 0.25 \left(\delta \mathbf{v}_{i+1}^{n} + 2 \delta \mathbf{v}_{i}^{n} + \delta \mathbf{v}_{i-1}^{n} \right) - \frac{\Delta t}{2\Delta x} \left(\mathbf{A}_{i+1}^{n} \cdot \delta \mathbf{v}_{i+1}^{n} - \mathbf{A}_{i-1}^{n} \cdot \delta \mathbf{v}_{i-1}^{n} \right)$$

$$+ \varsigma \frac{(\Delta t)^{2}}{2\Delta x^{2}} \left[\mathbf{D}_{i+0.5}^{n} \cdot \left(\delta \mathbf{v}_{i+1}^{n} - \delta \mathbf{v}_{i}^{n} \right) - \mathbf{D}_{i-0.5}^{n} \cdot \left(\delta \mathbf{v}_{i}^{n} - \delta \mathbf{v}_{i-1}^{n} \right) \right]$$

$$- \frac{\varkappa}{8} \left[\mathbb{B}_{i}^{n} \cdot \left(\delta \mathbf{v}_{i+1}^{n} - \delta \mathbf{v}_{i}^{n} \right) - \mathbb{A}_{i-1}^{n} \cdot \left(\delta \mathbf{v}_{i}^{n} - \delta \mathbf{v}_{i-1}^{n} \right) \right]. \tag{75}$$

Since the bounded operator $\mathbf{A} = \partial \mathbf{f}(\mathbf{u})/\partial \mathbf{u}$ in (22) is Fréchet-differentiable on the set $\Omega_{\mathbf{u}} \subset \mathbb{R}^M$, and its derivative is bounded on $\Omega_{\mathbf{u}}$, the operator \mathbf{A} is Lipschitz. In view of (69), we can write that $\|(\Delta t/\Delta x)(\mathbf{A}_{i-1}^n - \mathbf{A}_{i+1}^n) \cdot \delta \mathbf{v}_i^n\| \le C_a \Delta t \|\delta \mathbf{v}_i^n\|$, $C_a = const$. Thus, adding the term $(\Delta t/\Delta x)(\mathbf{A}_{i-1}^n - \mathbf{A}_{i+1}^n) \cdot \delta \mathbf{v}_i^n$ to the right-hand side of (75) has no effect on the stability of (75). To investigate the stability of (75) we rewrite it in the following form, where the term $0.5(\Delta t/\Delta x)(\mathbf{A}_{i-1}^n - \mathbf{A}_{i+1}^n) \cdot \delta \mathbf{v}_i^n$ is added.

$$\delta \mathbf{v}_{i}^{n+1} = \left(0.25\mathbf{I} - \mathbf{E}_{i,i-1}^{n}\right) \cdot \delta \mathbf{v}_{i-1}^{n}$$

$$+ \left(0.5\mathbf{I} + \mathbf{E}_{i,i-1}^{n} + \mathbf{E}_{i,i+1}^{n}\right) \cdot \delta \mathbf{v}_{i}^{n} + \left(0.25\mathbf{I} - \mathbf{E}_{i,i+1}^{n}\right) \cdot \delta \mathbf{v}_{i+1}^{n},$$

$$(76)$$

where

$$\mathbf{E}_{i,i-1}^{n} = \frac{\varkappa}{8} \mathbb{A}_{i-1}^{n} - \varsigma \frac{(\Delta t)^{2}}{2\Delta x^{2}} \mathbf{D}_{i-0.5}^{n} - \frac{\Delta t}{2\Delta x} \mathbf{A}_{i-1}^{n}, \tag{77}$$

$$\mathbf{E}_{i,i+1}^{n} = \frac{\varkappa}{8} \mathbb{B}_{i}^{n} - \varsigma \frac{\left(\Delta t\right)^{2}}{2\Delta x^{2}} \mathbf{D}_{i+0.5}^{n} + \frac{\Delta t}{2\Delta x} \mathbf{A}_{i+1}^{n}. \tag{78}$$

We write, in view of (23) and (42), that the spectrum $s\left(\mathbf{E}_{i,j}^n\right) \subseteq \left[L_{\min,i}^n, L_{\max,i}^n\right]$, where $j=i\pm 1$ and

$$L_{{\rm min},i}^n \ge -0.5\varsigma C_r^2 - 0.5C_r, \ L_{{\rm max},i}^n \le 0.5\varkappa \aleph - 0.25\varsigma C_r^2 + 0.5C_r. \eqno(79)$$

Hence, by virtue of [7, Theorem 2.10] we find that the scheme (76) will be stable if

$$\max_{\lambda \in \left[L_{\min,i}^n, L_{\max,i}^n\right]} (|0.25 - \lambda| + |0.25 + \lambda|) \leqslant 0.5, \quad \forall i, n.$$
 (80)

We obtain from (80) the following sufficient conditions for the stability of the variational scheme (71)

$$-0.5\zeta C_r^2 - 0.5C_r \ge -0.25,\tag{81}$$

$$0.5 \times \aleph - 0.25 \zeta C_r^2 + 0.5 C_r \le 0.25, \quad C_r = \frac{\Delta t \lambda_{\text{max}}}{\Delta x}.$$
 (82)

Let us note that the value of $L_{\max,i}^n$ in (79) can rarely be attained. It could be the case at the point (x_i, t_n) only if $\|\mathbb{B}_i^n\|_2 = 4\aleph$, $\|\mathbf{A}_{i+1}^n\|_2 = \lambda_{\max}$, and $\|\mathbf{A}_i^n\|_2 = 0$ (or, what is the same, $\|\mathbb{A}_{i-1}^n\|_2 = 4\aleph$, $\|\mathbf{A}_{i-1}^n\|_2 = \lambda_{\max}$, and $\|\mathbf{A}_i^n\|_2 = 0$). Moreover, considering the modification of the scheme (76), where $\mathbf{D}_{i-0.5}^n$ is replaced by $(\mathbf{A}_{i-1}^n)^2$, namely

$$\mathbf{E}_{i,i-1}^{n} = \frac{\varkappa}{8} \mathbb{A}_{i-1}^{n} - \varsigma \frac{(\Delta t)^{2}}{2\Delta x^{2}} (\mathbf{A}_{i-1}^{n})^{2} - \frac{\Delta t}{2\Delta x} \mathbf{A}_{i-1}^{n}, \tag{83}$$

$$\mathbf{E}_{i,i+1}^{n} = \frac{\varkappa}{8} \mathbb{B}_{i}^{n} - \varsigma \frac{(\Delta t)^{2}}{2\Delta x^{2}} (\mathbf{A}_{i+1}^{n})^{2} + \frac{\Delta t}{2\Delta x} \mathbf{A}_{i+1}^{n}, \tag{84}$$

we note (see Proof in [7, Theorem 2.10]) that the stability of the modified scheme, (76), (83)-(84), implies the stability of the scheme (76)-(78), since the operator $\bf A$ is Lipschitz. Hence, instead of (82), we will consider less rigid requirement:

$$0.5 \times \aleph - 0.5 \zeta C_r^2 + 0.5 C_r \le 0.25. \tag{85}$$

Let $\varsigma = 1$, then, by virtue of (81) and (85), we find the stability condition for the variational scheme (71):

$$C_r \le 0.5 \left(\sqrt{3} - 1\right), \quad \varkappa \aleph \le 2 - \sqrt{3}.$$
 (86)

Thus, in view of Theorem 8, the scheme MAC2, (64), will be stable if (86) will be valid.

Notice, if $\varsigma \to 0$ in (81) and (82), then we obtain the stability condition for the scheme MAC1, (50):

$$\varkappa \aleph + C_r \le 0.5, \quad C_r = \frac{\Delta t \lambda_{\text{max}}}{\Delta x}.$$
 (87)

It is significant that the stability conditions (81), (85), and hence (86), is found by virtue of [7, Theorem 2.10] and, therefore, it is assumed that $\mathbf{E}_{i,j}^n$ is Lipschitz, i.e. $\mathbf{E}_{i,i+1}^n = \mathbf{E}_{i,i-1}^n + O(\Delta x)$. For the sake of brevity, we associate the Landau symbol $O(\Delta x)$ with Lipschitz-continuity of a grid function whose increment can be estimated in norm by a constant times Δx for Δx small enough. Note that $\mathbf{E}_{i,j}^n$ is Lipschitz if \mathbf{A}_i^n , \mathbf{A}_i^n , and \mathbf{B}_i^n are Lipschitz. It is easy to see that \mathbf{A}_i^n will be Lipschitz if $\mathbf{u}_{\epsilon}^n(x)$ (the solution to (65) at $t=t_n$) will be Lipschitz-continuous. Actually, since the bounded operator $\mathbf{A}(\mathbf{u}_{\epsilon}) = (\partial \mathbf{f}/\partial \mathbf{u})|_{\mathbf{u}=\mathbf{u}_{\epsilon}}$ is Fréchet-differentiable on the set $\Omega_{\mathbf{u}} \subset \mathbb{R}^M$, and its derivative is bounded on $\Omega_{\mathbf{u}}$, we write: $\|\mathbf{A}_{i+1}^n - \mathbf{A}_i^n\| \leq C_A \|\mathbf{u}_{\epsilon,i+1}^n - \mathbf{u}_{\epsilon,i}^n\| \leq C_A C_u \Delta x$, where C_A , C_u

= const. We may guarantee that \mathbb{A}_i^n and \mathbb{B}_i^n will be Lipschitz if $\mathbf{u}_{\epsilon}^n(x)$ will be differentiable. Actually, we may write, by virtue of (28), that $\mathbb{A}_i = \mathbf{I} + O(\Delta x)$, and $\mathbb{B}_i = \mathbf{I} + O(\Delta x)$. Hence, Lipschitz-continuity of \mathbb{A}_i^n and \mathbb{B}_i^n is obvious. However, if $\mathbf{u}_{\epsilon}^n(x)$ is only Lipschitz-continuous, then we can not guarantee that \mathbb{A}_i^n and \mathbb{B}_i^n are also Lipschitz-continuous. In such a case a generalization of [7, Theorem 2.10] can be useful, namely the case when the scheme coefficients depend on several matrices. Then we obtain that (76), (83)-(84) will be stable if

$$\max_{\lambda_1, \lambda_2 \in \left[L_{\min,i}^n, L_{\max,i}^n\right]} (|0.25 - \lambda_1| + |0.25 - \lambda_2| + |0.5 + \lambda_1 + \lambda_2|) \leqslant 1, \quad \forall i, n. \quad (88)$$

By virtue of (88), we find that the scheme MAC2, (64), will be stable under the same conditions (81), (85), and, hence, under (86).

Let us note that the stability conditions, (86), are, in general, apt to be unduly rigid. Assuming that $\aleph = \aleph(x_i, t_n)$ as well as $C_r = C_r(t_n)$, and using Weyl's inequalities (see, e.g., [22, Chap. 3]), we can find less rigid stability conditions than (86). Let $\lambda_i^{n,k}$ denote the k-th singular value of \mathbf{A}_i^n , $\lambda_{\max,i}^n = \max_k \lambda_i^{n,k}$, $\lambda_{\max}^n = \max_i \lambda_{\max,i}^n$, and let $\alpha_i^{n,k}$ denote the k-th singular value of \mathbb{A}_i^n , $\alpha_{\min,i}^n = \min_k \alpha_i^{n,k}$. Then, given the pre-assigned $\aleph = \widehat{\aleph} \equiv const$ and $C_r = \widehat{C}_r \equiv const \leq 0.5$, we find the sought-after stability conditions at $t = t_n$, employing the following two-step procedure.

Algorithm 10

Step 1. Given $C_r = \widehat{C}_r \leq 0.5$ and $\tau^n = \Delta x \widehat{C}_r / \lambda_{\max}^n$, we, by virtue of Weyl's inequalities [22, Chap. 3], estimate $L_{\max,i}^n$:

$$L_{\max,i}^n \le 0.5 \varkappa \aleph_i^n - 0.5 \left(\frac{\tau^n}{\Delta x}\right)^2 \left(\lambda_{\max,i}^n\right)^2 + 0.5 \frac{\tau^n}{\Delta x} \lambda_{\max,i}^n. \tag{89}$$

In view of (88) and (89), we find

$$\aleph_i^n = \min\left[\widehat{\aleph}, \left(\frac{\tau^n}{\Delta x}\right)^2 \frac{\left(\lambda_{\max,i}^n\right)^2}{\varkappa} - \frac{\tau^n}{\varkappa \Delta x} \lambda_{\max,i}^n + \frac{1}{2\varkappa}\right]. \tag{90}$$

Step 2.Given \aleph_i^n , (90), we find the derivative \mathbf{d}_i^n using (43) and, by virtue of (28), the diagonal matrix \mathbb{A}_i^n , and hence $\alpha_{\min,i}^n$. We, by virtue of Weyl's inequalities [22, Chap. 3], estimate $L_{\min,i}^n$:

$$L_{\min,i}^n \ge \frac{\varkappa \alpha_{\min,i}^n}{8} - 0.5 \left(\frac{\tau^n}{\Delta x}\right)^2 \left(\lambda_{\max,i}^n\right)^2 - 0.5 \frac{\tau^n}{\Delta x} \lambda_{\max,i}^n. \tag{91}$$

Let τ_i^n be used instead of τ^n in (91). In view of (88) and (91), the following inequality must be valid.

$$\frac{\varkappa \alpha_{\min,i}^n}{4} - \left(\frac{\tau_i^n}{\Delta x}\right)^2 \left(\lambda_{\max,i}^n\right)^2 - \frac{\tau_i^n}{\Delta x} \lambda_{\max,i}^n \ge -0.5. \tag{92}$$

We obtain from (92):

$$\tau_i^n \frac{\lambda_{\max,i}^n}{\Delta x} \le \frac{\sqrt{3 + \varkappa \alpha_{\min,i}^n} - 1}{2}.$$
 (93)

Thus, given τ^n (= $\Delta x \hat{C}_r / \lambda_{\text{max}}^n$) and τ_i^n , (93), we find the time increment, Δt^n , at $t = t_n$:

$$\Delta t^n = \min\left(\tau^n, \min_i \tau_i^n\right) \implies C_r^n = \frac{\Delta t^n \lambda_{\max}^n}{\Delta x}.$$
 (94)

4 Operator splitting schemes

By virtue of the operator-splitting idea [6], [13], [19], [32] (see also LOS in [48]), the following chain of equations corresponds to the problem (1)

$$\frac{1}{2}\frac{\partial \mathbf{U}}{\partial t} = \frac{1}{\tau}\mathbf{q}(\mathbf{U}), \quad t_n < t \le t_{n+0.5}, \quad \mathbf{U}(\mathbf{x}, t_n) = \mathbf{U}^n(\mathbf{x}),$$
 (95)

$$\frac{1}{2}\frac{\partial \mathbf{U}}{\partial t} + \sum_{j=1}^{N} \frac{\partial}{\partial x_{j}} \mathbf{f}_{j} \left(\mathbf{U} \right) = 0, \ t_{n+0.5} < t \le t_{n+1}, \ \mathbf{U} \left(\mathbf{x}, t_{n+0.5} \right) = \mathbf{U}^{n+0.5} \left(\mathbf{x} \right), \ (96)$$

where $\mathbf{U}^n(\mathbf{x})$ denotes the solution to (96) at $t = t_n$, $\mathbf{U}^{n+0.5}(\mathbf{x})$ denotes the solution to (95) at $t = t_{n+0.5}$. If a high-resolution method is used directly for the homogeneous conservation law (96), then it is natural to use a high-order scheme for (95). As applied to, in general, stiff ($\tau \ll 1$) System (1), the second order schemes can be constructed on the basis of operator-splitting techniques with ease if (95) will be approximated by an implicit scheme and (96) by an explicit one, see Proposition 4.2 in [7]. As an example, let us develop a central scheme for a 1-D version of (1). After operator-splitting, the 1-D equation can be represented by the chain of equations, namely (95) and

$$\frac{1}{2}\frac{\partial \mathbf{U}}{\partial t} + \frac{\partial}{\partial x}\mathbf{f}\left(\mathbf{U}\right) = 0, \ t_{n+0.5} < t \le t_{n+1}, \ \mathbf{U}\left(x, t_{n+0.5}\right) = \mathbf{U}^{n+0.5}\left(x\right). \tag{97}$$

Let us first consider the case when the following first-order implicit scheme be used for (95)

$$\mathbf{v}_i^{n+0.5} = \mathbf{v}_i^n + \frac{\Delta t}{\tau} \mathbf{q} \left(\mathbf{v}_i^{n+0.5} \right), \tag{98}$$

and a central scheme with nonstaggered grid cells will be used for (97). We rewrite the scheme MAC1, (50), to read

$$\mathbf{v}_{i}^{n+1} = 0.25 \left(\mathbf{v}_{i-1}^{n+0.5} + 2\mathbf{v}_{i}^{n+0.5} + \mathbf{v}_{i+1}^{n+0.5} \right) - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n+0.5} - \mathbf{d}_{i-1}^{n+0.5} \right)$$
$$- \frac{\Delta t}{2\Delta x} \left(\mathbf{f}_{i+1}^{n+0.5} - \mathbf{f}_{i-1}^{n+0.5} \right), \quad \varkappa = const, \ 0 \le \varkappa \le 1,$$
(99)

where $\mathbf{d}_i^{n+0.5}$ denotes the derivative of interpolant at $x = x_i$. It is clear that the scheme (99) approximates (97) with the accuracy $O(\Delta t + (\Delta x)^2)$, however, in view of Proposition 4.2 in [7], the scheme (98)-(99), taken as a whole, is of the second order approximation for the 1-D version of (1).

Let us develop another nonstaggered central scheme approximating a 1-D version of (1) with the accuracy $O((\Delta t)^2 + (\Delta x)^2)$ and such that its components (after operator splitting) will be of the second order. It can be done on the basis of the second order scheme (98), (99) with ease. Notice, adding to and subtracting from Equation (4.10) in [7] (rewritten for $t_n < t \le t_{n+1}$) the same quantity is equivalent to adding this quantity to (99) and subtracting it from (98). Let $0.125 (\Delta t)^2 (\partial^2 \mathbf{U} / \partial t^2)_i^{n+0.5}$ be this quantity, then we obtain the following scheme, instead of (98), (99),

$$\mathbf{v}_i^{n+0.5} = \mathbf{v}_i^n + \frac{\Delta t}{\tau} \mathbf{q}_i^{n+0.5} - \frac{(\Delta t)^2}{8} \left(\frac{\partial^2 \mathbf{U}}{\partial t^2} \right)_i^{n+0.5}, \tag{100}$$

$$\mathbf{v}_{i}^{n+1} = 0.25 \left(\mathbf{v}_{i-1}^{n+0.5} + 2\mathbf{v}_{i}^{n+0.5} + \mathbf{v}_{i+1}^{n+0.5} \right) - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n+0.5} - \mathbf{d}_{i-1}^{n+0.5} \right)$$
$$- \frac{\Delta t}{2\Delta x} \left(\mathbf{f}_{i+1}^{n+0.5} - \mathbf{f}_{i-1}^{n+0.5} \right) + \frac{(\Delta t)^{2}}{8} \left(\frac{\partial^{2} \mathbf{U}}{\partial t^{2}} \right)_{i}^{n+0.5} . \tag{101}$$

Thus, the scheme (100) as well as the scheme (101) are of the second order, and the scheme (100)-(101), taken as a whole, is of the second order as well. Using Taylor series expansion, and central differencing, we find

$$\mathbf{v}_{i}^{n+0.25} = \mathbf{v}_{i}^{n+0.5} - \frac{\Delta t}{4} \left(\frac{\partial \mathbf{U}}{\partial t} \right)_{i}^{n+0.5} +$$

$$\frac{1}{2} \left(\frac{\Delta t}{4} \right)^2 \left(\frac{\partial^2 \mathbf{U}}{\partial t^2} \right)_i^{n+0.5} + O\left((\Delta t)^3 \right), \tag{102}$$

$$\mathbf{v}_{i}^{n+0.5} = \mathbf{v}_{i}^{n} + \frac{\Delta t}{2} \left(\frac{\partial \mathbf{U}}{\partial t} \right)_{i}^{n+0.25} + O\left((\Delta t)^{3} \right). \tag{103}$$

We obtain, by virtue of (62), (97), that

$$\frac{\partial^2 \mathbf{U}}{\partial t^2} = -2\frac{\partial}{\partial t} \left(\frac{\partial \mathbf{f}}{\partial x} \right) = -2\frac{\partial}{\partial x} \left(\frac{\partial \mathbf{f}}{\partial t} \right) = 4\frac{\partial}{\partial x} \left(\mathbf{A}^2 \cdot \frac{\partial \mathbf{U}}{\partial x} \right), \tag{104}$$

where $\mathbf{A} = \partial \mathbf{f} / \partial \mathbf{U}$. Then

$$\left[\frac{\partial}{\partial x} \left(\mathbf{A}^2 \cdot \frac{\partial \mathbf{U}}{\partial x} \right) \right]_i^{n+0.5} = \frac{1}{\Delta x} \left[\mathbf{D}_{i+0.5}^{n+0.5} \cdot \frac{\mathbf{v}_{i+1}^{n+0.5} - \mathbf{v}_i^{n+0.5}}{\Delta x} \right]$$

$$- \mathbf{D}_{i-0.5}^{n+0.5} \cdot \frac{\mathbf{v}_{i}^{n+0.5} - \mathbf{v}_{i-1}^{n+0.5}}{\Delta x} \right] + O\left((\Delta x)^{2}\right), \tag{105}$$

where $\mathbf{D}_{i+0.5}^{n+0.5} = 0.5 \left((\mathbf{A}_{i+1}^{n+0.5})^2 + (\mathbf{A}_{i}^{n+0.5})^2 \right)$. By virtue of (95), (102)-(105), we rewrite Scheme (100)-(101) to read

$$\mathbf{v}_{i}^{n+0.25} = \mathbf{v}_{i}^{n+0.5} - \frac{\Delta t}{4\tau} \left(\mathbf{q}_{i}^{n+0.25} + \mathbf{q}_{i}^{n+0.5} \right), \tag{106}$$

$$\mathbf{v}_{i}^{n+0.5} = \mathbf{v}_{i}^{n} + \frac{\Delta t}{\tau} \mathbf{q}_{i}^{n+0.25}, \tag{107}$$

$$\mathbf{v}_{i}^{n+1} = 0.25 \left(\mathbf{v}_{i-1}^{n+0.5} + 2\mathbf{v}_{i}^{n+0.5} + \mathbf{v}_{i+1}^{n+0.5} \right) - \varkappa \frac{\Delta x}{8} \left(\mathbf{d}_{i+1}^{n+0.5} - \mathbf{d}_{i-1}^{n+0.5} \right)$$

$$+\frac{\varsigma\left(\Delta t\right)^{2}}{2\left(\Delta x\right)^{2}}\left[\mathbf{D}_{i+0.5}^{n+0.5}\cdot\left(\mathbf{v}_{i+1}^{n+0.5}-\mathbf{v}_{i}^{n+0.5}\right)-\mathbf{D}_{i-0.5}^{n+0.5}\cdot\left(\mathbf{v}_{i}^{n+0.5}-\mathbf{v}_{i-1}^{n+0.5}\right)\right]$$

$$-\frac{\Delta t}{2\Delta x} \left(\mathbf{f}_{i+1}^{n+0.5} - \mathbf{f}_{i-1}^{n+0.5} \right). \tag{108}$$

If $\varsigma = 0$, then (108) coincides with (99), being $O(\Delta t + (\Delta x)^2)$ accurate. If $\varkappa = 1$ and $\varsigma = 1$, then Scheme (106)-(108) approximates the 1-D version of (1) with the accuracy $O((\Delta t)^2 + (\Delta x)^2)$. The scheme (108) coincides, in fact, with the scheme MAC2, (64), and hence the scheme (108) will be stable if (86) will be valid. Notice, the less rigid conditions for the stability of the scheme (108) can be found by means of Algorithm 10.

Let us note that in practice (e.g., [32], [48]) the operator-splitting techniques find a wide range of application in designing economical schemes for Eq. (96) in the domains of complicated geometry. The resulting method, in general, will be only first-order accurate in time because of the splitting [32], [48]. Thus, in line with established practice we will replace the multidimensional Eq. (96) by the chain of the one-dimensional equations:

$$\frac{1}{2N}\frac{\partial \mathbf{U}_j}{\partial t} + \frac{\partial}{\partial x_j}\mathbf{f}_j(\mathbf{U}_j) = 0, \ t_{n+0.5+(j-1)/(2N)} < t \le t_{n+0.5+j/(2N)}, \quad (109)$$

where $\mathbf{U}_{j}\left(x,t_{n+0.5+(j-1)/(2N)}\right) = \mathbf{U}_{j-1}\left(x,t_{n+0.5+(j-1)/(2N)}\right), \ j=1,2,\ldots N,$ $\mathbf{U}_{0}\left(x,t_{n+0.5}\right)$ denotes the solution to (95) at $t=t_{n+0.5}$. Eq. (95) will be approximated by the first-order implicit scheme, (98), or a second-order implicit Runge-Kutta scheme. In particular, it can be used (106)-(107) or the following implicit Runge-Kutta scheme

$$\mathbf{v}_{i}^{n+0.25} = \mathbf{v}_{i}^{n+0.5} - \frac{\Delta t}{2\tau} \mathbf{q}_{i}^{n+0.5}, \ \mathbf{v}_{i}^{n+0.5} = \mathbf{v}_{i}^{n} + \frac{\Delta t}{\tau} \mathbf{q}_{i}^{n+0.25},$$
(110)

since these schemes possess a discrete analogy to the continuous asymptotic limit (see, e.g., [24], [37], [41]). Let us note that the scheme (110) is an implicit analogue to the well known Runge-Kutta scheme, which is referred, originally due to Runge, as modified Euler method [20]:

$$\mathbf{v}_{i}^{n+0.25} = \mathbf{v}_{i}^{n} + \frac{\Delta t}{2\tau} \mathbf{q}_{i}^{n}, \ \mathbf{v}_{i}^{n+0.5} = \mathbf{v}_{i}^{n} + \frac{\Delta t}{\tau} \mathbf{q}_{i}^{n+0.25}.$$
(111)

5 Examples

In this section, we are mainly concerned with verification of the second order central schemes COS2, (59) and MAC2, (64), as well as the splitting schemes (106)-(108), and (109).

5.1 Scalar non-linear equation

As the first stage in the verification, we will focus on the following scalar 1-D version of the problem (1):

$$\frac{\partial u}{\partial t} + \frac{\partial}{\partial x} f(u) = 0, \quad x \in \mathbb{R}, \ 0 < t \le T_{\text{max}}; \quad u(x,t)|_{t=0} = u^{0}(x). \tag{112}$$

Let us first compare the schemes MAC2, (64), and COS2, (59), and demonstrate that these schemes are of the second order. We consider the linear transport equation, i.e. (112) with $f(u) \equiv u$, subject to the initial data: $u^0(x) = \sin(\pi x)$. The numerical solutions were computed under $\varkappa = \varsigma = 1$. The scheme COS2, (59), will be stable, in view of (67), (68), if we take the CFL number $Cr = \sqrt{3} - 1$, and $\aleph = 2 - \sqrt{3}$. The scheme MAC2, (64), will be stable, in view of (86), under $Cr = 0.5(\sqrt{3} - 1)$ and $\aleph = 2 - \sqrt{3}$. We will also verify the scheme COS1, (45), under Cr = 0.5 and $\aleph = 0.25$, as well as the scheme MAC1, (50), under Cr = 0.25 and $\aleph = 0.25$. Notice, in the case of schemes COS2 and COS1 the CFL numbers are two times higher than in the case of schemes MAC2 and MAC1, respectively. The reason is that the staggered schemes, COS2 and COS1, are solved twice during the time increment, Δt , in contrast to the non-staggered schemes MAC2 and MAC1. L_1 errors, at t = 10, $0 \le x \le 2$, versus the number of nodes are depicted in Table 1.

Table 1: L_1 errors versus the number of nodes (N)

ı	N	1280	640	320	160	80
				9.8×10^{-3}		
				2.5×10^{-4}		
				1.0×10^{-2}		
	MAC2	5.6×10^{-5}	2.4×10^{-4}	9.8×10^{-4}	3.9×10^{-3}	1.5×10^{-2}

Note (Table 1) that L_1 errors in the case of the scheme COS2 are approximately coincide with those appeared in the case of the scheme MAC2 under the double number of nodes. It is associated with the fact that in the case of the staggered scheme COS2, (59), the space increment is, actually, two times less than in the case of the nonstaggered scheme MAC2, (64).

Now we will solve the inviscid Burgers equation (i.e. $f(u) \equiv u^2/2$) with the following initial condition

$$u(x,0) = \begin{cases} u_0, & x \in (h_L, h_R) \\ 0, & x \notin (h_L, h_R) \end{cases}, \quad h_R > h_L, \ u_0 = const \neq 0.$$
 (113)

The exact solution to (112), (113) is given by

$$u(x,t) = \begin{cases} u_1(x,t), & 0 < t \le T \\ u_2(x,t), & t > T \end{cases},$$
 (114)

where $T = 2S/u_0, S = h_R - h_L,$

$$u_1(x,t) = \begin{cases} \frac{x - h_L}{b - h_L} u_0, & h_L < x \le b, \ b = u_0 t + h_L \\ u_0, & b < x \le 0.5 u_0 t + h_R \\ 0, & x \le h_L \ or \ x > 0.5 u_0 t + h_R \end{cases} , \tag{115}$$

$$u_2(x,t) = \begin{cases} \frac{2S(x-h_L)}{(L-h_L)^2} u_0, & h_L < x \le L \\ 0, & x \le h_L \text{ or } x > L \end{cases},$$
(116)

$$L = 2\sqrt{S^2 + 0.5u_0S(t - T)} + h_L. \tag{117}$$

First, it will be used the first order in time schemes MAC1, (50), and COS1, (45). The numerical solutions were computed on a uniform grid with spatial increments of $\Delta x = 0.01$, the velocity $u_0 = 1$ in (113), $h_L = 0.2$, $h_R = 1$. In view of the stability condition, (87), for the scheme MAC1, we take the CFL number Cr = 0.25, the monotonicity parameter $\aleph = 0.25$, and the parameter $\varkappa = 1$. Having regard to the stability condition (67) under $\varsigma = 0$, we verify the scheme COS1, (45), under Cr = 0.5, $\aleph = 0.25$, and $\varkappa = 1$. The results of simulations are depicted with the exact solution in Figure 2.

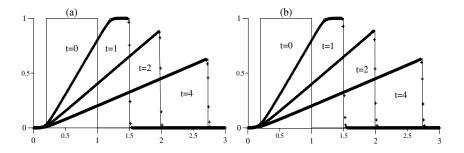


Figure 2: Inviscid Burgers equation. The schemes COS1 (a) at given $C_r = 0.5$ and MAC1 (b) at given $C_r = 0.25$ versus the analytical solution. Crosses: numerical solution; Solid line: analytical solution and initial data. $\Delta x = 0.01$, $\aleph = 0.25$.

We note (Figure 2) that the schemes COS1, (45), and MAC1, (50), exhibit similar results in spite of the fact that in the case of the scheme COS1 the space increments are, virtually, two times less than those in the case of MAC1. It should also be mentioned the absence of spurious oscillations in the numerical solutions (see Figure 2). However, if we take $\aleph = 0.5$, then the spurious oscillations can be produced by the scheme COS1 (see [8, p. 2804]). Using Algorithm 10 at given $\widehat{\aleph} = 0.5$, in the case of scheme MAC1 we obtain very slight but

significant spurious oscillations, which indicate that the boundary maximum principle is violated. For instance, the maximum positive value of the dependent variable, v, at t=1 is 0.4% above the maximum value of v at the boundary t=0. These spurious oscillations can be eradicated by decreasing the parameter \varkappa . Particularly, the spurious oscillations disappear if $\varkappa=0.7$, however, this introduces an additional numerical smearing. The results of simulations are not depicted here.

To test the schemes MAC2, (64), and COS2, (59), the inviscid Burgers equation was solved under the initial condition (113). The numerical solutions were computed under the same values of parameters as in the case of the schemes MAC1 and COS1, but C_r and \aleph . In view of the stability condition, (86), for the scheme MAC2, we take the CFL number $Cr = 0.5(\sqrt{3} - 1)$. It is also used Algorithm 10 at given $\widehat{\aleph} = 0.5$. The scheme COS2 will be stable, in view of (67), (68), if we take $Cr = \sqrt{3} - 1$ and $\aleph = 0.5$. The results of simulation are depicted with the exact solution in Figure 3.

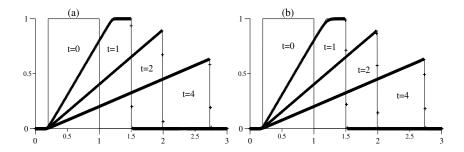


Figure 3: Inviscid Burgers equation. The schemes COS2 (a) at given $C_r = \sqrt{3} - 1$ and MAC2 (b) at given $C_r = 0.5 (\sqrt{3} - 1)$ versus the analytical solution. Crosses: numerical solution; Solid line: analytical solution and initial data. $\Delta x = 0.01$, $\aleph = \widehat{\aleph} = 0.5$.

We note (Figure 3) that the schemes COS2 and MAC2 exhibit a typical second-order nature without any spurious oscillations. The scheme MAC2 exhibits more numerical smearing than the scheme COS2. Such a phenomenon has, at least, two reasons. First, in the case of the staggered scheme COS2 the space increment is, in fact, two times less than in the case of the nonstaggered scheme MAC2. Second place, using Algorithm 10 at given $\hat{\aleph}=0.5$ we, in general, obtain $\aleph_i^n < \hat{\aleph}$ at the grid node (x_i,t_n) , resulting in an additional numerical smearing.

5.2 Hyperbolic conservation laws with relaxation

Let us consider the model system of hyperbolic conservation laws with relaxation developed in [45]:

$$\frac{\partial w}{\partial t} + \frac{\partial}{\partial x} \left(\frac{1}{2} u^2 + aw \right) = 0, \tag{118}$$

$$\frac{\partial z}{\partial t} + \frac{\partial}{\partial x} az = \frac{1}{\tau} Q(w, z), \tag{119}$$

where

$$Q(w,z) = z - m(u - u_0), \quad u = w - q_0 z, \tag{120}$$

 τ denotes the relaxation time of the system, q_0 , m, a, and u_0 are constants. The Jacobian, \mathbf{A} , can be written in the form

$$\mathbf{A} = \left\{ \begin{array}{cc} w - q_0 z + a & -q_0 \left(w - q_0 z \right) \\ 0 & a \end{array} \right\}. \tag{121}$$

The system (118)-(119) has the following frozen [45] characteristic speeds $\lambda_1 = a$, $\lambda_2 = u + a$. The equilibrium equation for (118)-(119) is

$$\frac{\partial w}{\partial t} + \frac{\partial}{\partial x} \left(\frac{1}{2} u_*^2 + aw \right) = 0, \tag{122}$$

where

$$u_* = w - q_0 z_*, \quad z_* = \frac{m}{1 + mq_0} (w - u_0).$$
 (123)

The equilibrium characteristic speed λ_* can be written in the form

$$\lambda_*(w) = \frac{u_*(w)}{1 + mq_0} + a.$$
 (124)

Pember's rarefaction test problem is to find the solution $\{w, z\}$ to (118)-(119), and hence the function u = u(x, t), under $\tau \to 0$, and where

$$\{w, z\} = \begin{cases} \{w_L, z_* (w_L)\}, & x < x_0 \\ \{w_R, z_* (w_R)\}, & x > x_0 \end{cases},$$
(125)

$$0 < u_L = w_L - q_0 z_* (w_L) < u_R = w_R - q_0 z_* (w_R).$$
(126)

The analytical solution of this problem can be found in [45]. The parameters of the model system are assumed as follows: $q_0 = -1$, m = -1, $u_0 = 3$, a = 1, $\tau = 10^{-8}$. The initial conditions of the rarefaction problem are defined by

$$u_L = 2, \implies z_L = m(u_L - u_0) = 1, \ w_L = u_L + q_0 z_L = 1,$$
 (127)

$$u_R = 3, \implies z_R = m(u_R - u_0) = 0, \ w_R = u_R + q_0 z_R = 3.$$
 (128)

The position of the initial discontinuity, x_0 , is set according to the value of a so that the solutions of all the rarefaction problems are identical [45]. Let a position, x_R^t , of leading edge or a position, x_L^t , of trailing edge of the rarefaction be known (e.g., $x_R^t = 0.85$, $x_L^t = 0.7$ in [45]), then

$$x_0 = x_R^t - \left(\frac{u_R}{1 + mq_0} + a\right)t = x_L^t - \left(\frac{u_L}{1 + mq_0} + a\right)t.$$
 (129)

At t = 0.3, under (127)-(128) we have [45]

$$u = \begin{cases} 2, & x \le 0.7\\ 2 + \frac{x - 0.7}{0.85 - 0.7}, & 0.7 < x < 0.85\\ 3, & x \ge 0.85 \end{cases}$$
 (130)

The results of simulations, based upon the schemes MAC2, (64), and COS2, (59), together with (110), under different values of a grid spacing ($\Delta x = 10^{-3}$, 5×10^{-4}) are depicted in Figure 4. In view of the stability condition, (86), for the scheme MAC2, we take the CFL number $Cr = 0.5(\sqrt{3} - 1)$. It is also used Algorithm 10 at given $\hat{\aleph} = 0.5$. The scheme COS2 will be stable, in view of (67), (68), if we take $Cr = \sqrt{3} - 1$ and $\aleph = 0.5$.

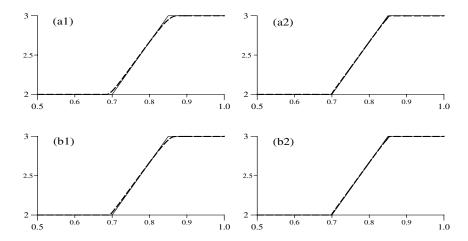


Figure 4: Pember's rarefaction test problem. The schemes COS2 (a1, a2) at given $C_r = \sqrt{3} - 1$ and MAC2 (b1, b2) at given $C_r = 0.5 (\sqrt{3} - 1)$ versus the analytical solution for u. Dashed line: numerical solution; Solid line: analytical solution. Time t = 0.3, $\aleph = \widehat{\aleph} = 0.5$, (a1, b1): $\Delta x = 0.001$, (a2, b2): $\Delta x = 0.0005$.

One can clearly see (Figure 4) that the schemes MAC2 together with the Runge-Kutta scheme (110) as well as COS2 together with (110) are free from spurious oscillations. Recall that in the case of the staggered scheme COS2 the space increment is, in fact, two times less than in the case of the nonstaggered scheme MAC2. Moreover, using Algorithm 10 in the case of scheme MAC2, we, in general, obtain $\aleph_i^n < \widehat{\aleph}$ at the grid node (x_i, t_n) , resulting in an additional numerical smearing. Nevertheless, the scheme MAC2 can exhibit (Figure 4) even less numerical viscosity than the scheme COS2.

5.3 1-D Euler equation of gas dynamics

In this subsection we apply the second order schemes COS2, (59), and MAC2, (64), to the Euler equations of gamma-law gas:

$$\frac{\partial \mathbf{u}(x,t)}{\partial t} + \frac{\partial}{\partial x} \mathbf{F}(\mathbf{u}) = 0, \quad x \in \mathbb{R}, \ t > 0; \quad \mathbf{u}(x,0) = \mathbf{u}^{0}(x), \tag{131}$$

$$\mathbf{u} \equiv \{u_1, u_2, u_3\}^T = \{\rho, \rho v, e\}^T, \quad \mathbf{F}(\mathbf{u}) = \{\rho v, \rho v^2 + p, (e+p)v\}^T, \quad (132)$$

$$e = \frac{p}{\gamma - 1} + \frac{1}{2}\rho v^2, \quad \gamma = const, \tag{133}$$

where ρ , v, p, e denote the density, velocity, pressure, and total energy respectively. We consider the Riemann problem subject to Riemann initial data

$$\mathbf{u}^{0}(x) = \begin{cases} \mathbf{u}_{L} & x < x_{0} \\ \mathbf{u}_{R} & x > x_{0} \end{cases}, \quad \mathbf{u}_{L}, \mathbf{u}_{R} = const.$$
 (134)

The analytic solution to the Riemann problem can be found in [32, Sec. 14].

We solve the shock tube problem (see, e.g., [5], [23], [32], [33]) with Sod's initial data:

$$\mathbf{u}_{L} = \left\{ \begin{array}{c} 1\\0\\2.5 \end{array} \right\}, \quad \mathbf{u}_{R} = \left\{ \begin{array}{c} 0.125\\0\\0.25 \end{array} \right\}. \tag{135}$$

Following Balaguer and Conde [5] as well as Liu and Tadmor [33] we assume that the computational domain is $0 \le x \le 1$; the point x_0 is located at the middle of the interval [0, 1], i.e. $x_0 = 0.5$; the equations (131) are integrated up to t = 0.16 on a spatial grid with 200 nodes as in [5] and in [33]. Aiming to compare the schemes COS2 and MAC2, we, in view of the stability conditions (67), (68), and (86), take the CFL number $Cr = \sqrt{3} - 1$ (the scheme COS2), and $Cr = 0.5(\sqrt{3} - 1)$ (the scheme MAC2). In both cases we take $\varkappa \aleph = 2 - \sqrt{3}$, where $\varkappa = 0.94$. The reason why the parameter \varkappa has to be less than unity is discussed in Section 3.2 (see also [8, Sec. 4]). The results of simulations are depicted in Figure 5 (left column) and Figure 6 (left column).

The results using the scheme COS2 (Figure 5, left column) are not worse in comparison to the corresponding third-order central results of [33, p. 418] as well as to the results obtained by the fourth-order non-oscillatory scheme in [5, p. 472]. The simulations in [5] and [33] were done under $\Delta t = 0.1\Delta x$ (i.e. $0.13 \lesssim Cr \lesssim 0.22$). The fourth-order scheme [5, p. 472] gives a better resolution but, in contrast to the scheme COS2 (Figure 5, left column), can produce spurious oscillations. Let us note that the second-order scheme COS2 can give results analogous to those obtained by the fourth-order scheme [5, p. 472]. Namely, the higher will be the values of \varkappa and \aleph , the better will be the resolution (see, e.g., Figure 5, right column). However, if $\varkappa=1$, then the scheme COS2 will be oscillatory [8, p. 2816]. It is apparent that \varkappa should be reasonably less than unity to suppress spurious oscillations, while it should be as close to unity as possible to lose less in the order of accuracy of the scheme [8,

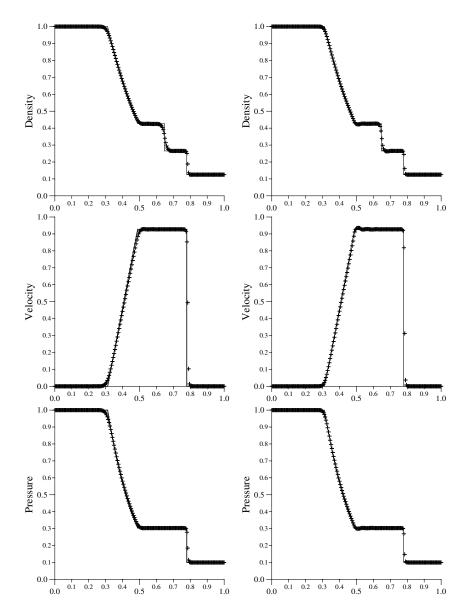


Figure 5: Sod's problem. Time t=0.16. The scheme COS2 versus the analytical solution. $\varkappa\aleph=2-\sqrt{3},\ \varkappa=0.94$ (left column) and $\aleph=0.35,\ \varkappa=1$ (right column). Crosses: numerical solution; Solid line: analytical solution. $C_r=\sqrt{3}-1,\ \Delta x=0.005$

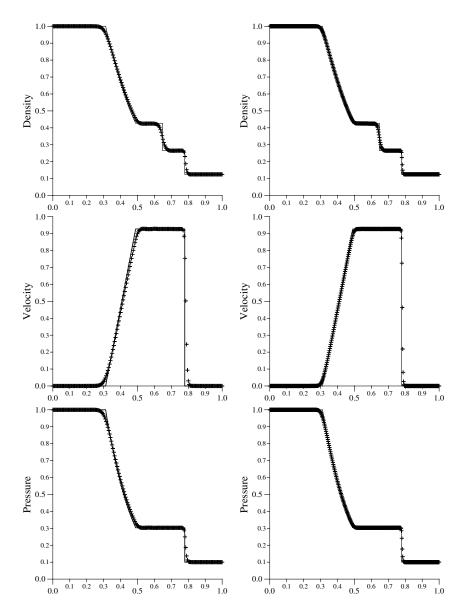


Figure 6: Sod's problem. The scheme MAC2 versus the analytical solution. Time t=0.16. $C_r=0.5(\sqrt{3}-1)$, $\varkappa\aleph=2-\sqrt{3}$, $\varkappa=0.94$. Crosses: numerical solution; Solid line: analytical solution. $\Delta x=0.005$ (left column), $\Delta x=0.0025$ (right column).

p. 2816]. The results of using the scheme MAC2 (Figure 6, left column) are not worse in comparison to the corresponding second-order central results of [33, p. 418], in spite of the fact in the case of staggered scheme the space increment is, in effect, two times less than in the case of nonstaggered scheme. For the sake of illustration, the results using the scheme MAC2 under $\Delta x = 0.0025$ are depicted in Figure 6 (right column). It is easily seen that the scheme MAC2 gives a better resolution than the scheme COS2 (Figure 5, left column).

5.4 3-D axial symmetric gas dynamics

We consider an adiabatic expansion of a vapor-plasma plume, produced by a laser beam on a target plane surface, into vacuum [3], i.e., the so called Anisimov's problem. It is assumed [3] that the expansion of the plasma plume may be described by the ideal gas dynamic equations with a constant adiabatic exponent γ . Let the target surface be perpendicular to the axis z and located at z=0. Taking into account the symmetry of the plume with respect to the axis z, the gas-dynamic equations can be written as follows.

$$\frac{\partial \rho}{\partial t} + \frac{1}{r} \frac{\partial (r\rho v_r)}{\partial r} + \frac{\partial (\rho v_z)}{\partial z} = 0, \tag{136}$$

$$\frac{\partial}{\partial t} (\rho v_r) + \frac{1}{r} \frac{\partial}{\partial r} \left[r \rho (v_r)^2 \right] + \frac{\partial}{\partial z} (\rho v_z v_r) + \frac{\partial p}{\partial r} = 0, \tag{137}$$

$$\frac{\partial}{\partial t} (\rho v_z) + \frac{\partial}{\partial z} \left[\rho (v_z)^2 \right] + \frac{1}{r} \frac{\partial}{\partial r} (r \rho v_z v_r) + \frac{\partial p}{\partial z} = 0, \tag{138}$$

$$\frac{\partial \rho E}{\partial t} + \frac{1}{r} \frac{\partial}{\partial r} \left[r v_r \left(\rho E + p \right) \right] + \frac{\partial}{\partial z} \left[v_z \left(\rho E + p \right) \right] = 0. \tag{139}$$

$$\rho E = \frac{p}{\gamma - 1} + 0.5\rho v^2, \ v^2 = v_r^2 + v_z^2, \ \gamma = const. \tag{140}$$

Let R_0 and Z_0 will be the initial radius of the plume at z=0 and the initial height of the plume in z-direction, respectively. It is also used in [3] the mass M_P of the plume and its initial energy E_P as the input data. We will use the following values as the reference quantities: $l_* = R_0$, $v_* = \sqrt{(5\gamma - 3)} E_P/M_P$, $t_* = l_*/v_*$, $\rho_* = M_P/(R_0^2 Z_0)$, $p_* = \rho_* v_*^2$. It is assumed that the system (136)-(139) is written in the non-dimensional variables, since the dimensionless form of the gas dynamic equations coincides with (136)-(139), up to notations. The pseudo-analytic solution to (136)-(139) can be written in the following form [3]:

$$\rho(r,z,t) = I_1^{-1}(\gamma) \frac{\sigma}{\xi^2 \eta} \left(1 - \frac{r^2}{\xi^2} - \frac{z^2}{\eta^2} \right)^{1/(\gamma - 1)}, \ \xi = \xi(t), \ \eta = \eta(t), \quad (141)$$

$$p(r,z,t) = \frac{I_2^{-1}(\gamma)}{5\gamma - 3} \left(\frac{\sigma}{\xi^2 \eta}\right)^{\gamma} \left(1 - \frac{r^2}{\xi^2} - \frac{z^2}{\eta^2}\right)^{\gamma/(\gamma - 1)}, \quad \sigma = \frac{Z_0}{R_0}, \tag{142}$$

$$v_r(r,z,t) = r\frac{\dot{\xi}}{\xi}, \ v_z(r,z,t) = z\frac{\dot{\eta}}{\eta}, \quad \dot{\xi} \equiv \frac{d\xi}{dt}, \ \dot{\eta} \equiv \frac{d\eta}{dt},$$
 (143)

$$\xi \ddot{\xi} = \eta \ddot{\eta} = \left(\frac{\sigma}{\xi^2 \eta}\right)^{\gamma - 1}, \quad \xi(0) = 1, \ \eta(0) = \sigma, \ \dot{\xi}(0) = \dot{\eta}(0) = 0,$$
 (144)

where

$$I_{1}\left(\gamma\right) = \frac{\pi^{3/2}}{2} \Gamma\left(\frac{\gamma}{\gamma - 1}\right) / \Gamma\left(\frac{\gamma}{\gamma - 1} + \frac{3}{2}\right),\tag{145}$$

$$I_{2}\left(\gamma\right) = \frac{\pi^{3/2}}{2\left(\gamma - 1\right)} \Gamma\left(\frac{\gamma}{\gamma - 1} + 1\right) / \Gamma\left(\frac{\gamma}{\gamma - 1} + \frac{5}{2}\right),\tag{146}$$

 $\Gamma(\cdot)$ is the Gamma-function. Out of the plume, i.e. provided $(r/\xi)^2 + (z/\eta)^2 > 1$, the values of density (ρ) and pressure (p) as well as the components of velocity (v_r, v_z) are equal to zero.

Thus, the problem, (136)-(139), is reduced to the system of ordinary differential equations (ODEs) (144). The system of ODEs, (144), is solved numerically by the Runge-Kutta method with the adaptive step-size control [46]. To check the accuracy of the calculations, the integral of energy [3] is used:

$$\dot{\xi}^2 + 0.5\dot{\eta}^2 + \frac{1}{\gamma - 1} \left(\frac{\sigma}{\xi^2 \eta}\right)^{\gamma - 1} = \frac{1}{\gamma - 1}.$$
 (147)

Aiming to use the schemes COS2, (59), and MAC2, (64), for solving the system (136)-(139), we rewrite (136)-(139) to read

$$\frac{\partial \rho}{\partial t} + \frac{\partial (\rho v_r)}{\partial r} + \frac{\partial (\rho v_z)}{\partial z} = -\frac{\rho v_r}{r},\tag{148}$$

$$\frac{\partial}{\partial t} (\rho v_r) + \frac{\partial}{\partial r} \left[\rho (v_r)^2 + p \right] + \frac{\partial}{\partial z} (\rho v_z v_r) = -\frac{\rho (v_r)^2}{r}, \tag{149}$$

$$\frac{\partial}{\partial t} (\rho v_z) + \frac{\partial}{\partial z} \left[\rho (v_z)^2 + p \right] + \frac{\partial}{\partial r} (\rho v_z v_r) = -\frac{\rho v_z v_r}{r}, \tag{150}$$

$$\frac{\partial \rho E}{\partial t} + \frac{\partial}{\partial r} \left[v_r \left(\rho E + p \right) \right] + \frac{\partial}{\partial z} \left[v_z \left(\rho E + p \right) \right] = -\frac{v_r \left(\rho E + p \right)}{r}. \tag{151}$$

The initial conditions are the following (in details, see (141)-(146) under t = 0):

$$v_r = v_z = 0, \ \rho = I_1^{-1}(\gamma) \left(1 - r^2 - (z/\sigma)^2\right)^{1/(\gamma - 1)}, \ p/\rho^{\gamma} = const.$$
 (152)

The boundary conditions are established on the basis of the reflection concept [2]. At r=0 we assume that the axis z is a reflection line. It prohibits any normal flux of mass through the boundary r=0, i.e. $v_r=0$. Moreover, it is assumed that the pressure (p), density (ρ) , and tangential velocity (v_z) are even functions of normal distance to the axis z while the normal velocity (v_r) is an odd function of r. It is also assumed that the plane z=0 is a reflection surface, i.e. the pressure (p), density (ρ) , and tangential velocity (v_r) are even functions of normal distance above the target surface while the normal velocity (v_z) is an odd function of z.

By virtue of the operator-splitting method [6], [13], [19], [32] (see also LOS in [48]), the system (148)-(151) may be approximated by the following chain of equations:

$$\frac{1}{4}\frac{\partial \check{\mathbf{u}}}{\partial t} + \frac{\partial}{\partial z}\check{\mathbf{f}}\left(\check{\mathbf{u}}\right) = 0, \quad t_n < t \le t_{n+0.25}, \ \check{\mathbf{u}}|_{t=t_n} = \hat{\mathbf{u}}^n, \tag{153}$$

$$\frac{1}{4}\frac{\partial \bar{\mathbf{u}}}{\partial t} + \frac{\partial}{\partial r}\bar{\mathbf{f}}(\bar{\mathbf{u}}) = 0, \ t_{n+0.25} < t \le t_{n+0.5}, \ \bar{\mathbf{u}}|_{t=t_{n+0.25}} = \check{\mathbf{u}}^{n+0.25},$$
(154)

$$\frac{1}{2}\frac{\partial \hat{\mathbf{u}}}{\partial t} = -\frac{1}{r}\mathbf{q}(\hat{\mathbf{u}}), \ t_{n+0.5} < t \le t_{n+1}, \ \hat{\mathbf{u}}|_{t=t_{n+0.5}} = \bar{\mathbf{u}}^{n+0.5},$$
 (155)

where
$$\mathbf{\check{u}} = \{\check{\rho}, \check{\rho}\check{v}_r, \check{\rho}\check{v}_z, \check{\rho}\check{E}\}^T$$
, $\mathbf{\check{f}}(\check{\mathbf{u}}) = \{\check{\rho}\check{v}_z, \check{\rho}\check{v}_z\check{v}_r, \check{\rho}(\check{v}_z)^2 + \check{p}, \check{v}_z(\check{\rho}\check{E} + \check{p})\}^T$, $\mathbf{\bar{u}} = \{\bar{\rho}, \bar{\rho}\bar{v}_r, \bar{\rho}\bar{v}_z, \bar{\rho}\bar{E}\}^T$, $\mathbf{\bar{f}}(\bar{\mathbf{u}}) = \{\bar{\rho}\bar{v}_r, \bar{\rho}(\bar{v}_r)^2 + \bar{p}, \bar{\rho}\bar{v}_z\bar{v}_r, \bar{v}_r(\bar{\rho}\bar{E} + \bar{p})\}^T$, $\mathbf{\hat{u}} = \{\hat{\rho}, \hat{\rho}\hat{v}_r, \hat{\rho}\hat{v}_z, \hat{\rho}\hat{E}\}^T$, $\mathbf{q}(\hat{\mathbf{u}}) = \{\hat{\rho}\hat{v}_r, \hat{\rho}(\hat{v}_r)^2, \hat{\rho}\hat{v}_z\hat{v}_r, \hat{v}_r(\hat{\rho}\hat{E} + \hat{p})\}^T$.

We will use the schemes COS2, (59), and MAC2, (64), fo solving Eqs. (153)-(154). To solve Eq. (155) it will be used the Runge-Kutta scheme (modified Euler method), (111). Let us note that every point on the axis r=0 is a singular point for Eq. (155). Assuming that all terms at left-hand side of Equation (148) are bounded values at a vicinity of r=0, we find that $v_r \to 0$ as $r \to 0$ and, hence, $\mathbf{q}(\hat{\mathbf{u}})|_{r=0}=0$. Thus

$$\lim_{r\to 0} \frac{\mathbf{q}\left(\hat{\mathbf{u}}\right)|_{r>0}}{r} = \lim_{r\to 0} \frac{\mathbf{q}\left(\hat{\mathbf{u}}\right)|_{r>0} - \mathbf{q}\left(\hat{\mathbf{u}}\right)|_{r=0}}{r} = \left. \frac{\partial \mathbf{q}\left(\hat{\mathbf{u}}\right)}{\partial r} \right|_{r=0}.$$
 (156)

Thus, we have at r=0:

$$\frac{1}{2}\frac{\partial \hat{\mathbf{u}}}{\partial t} = -\frac{\partial \mathbf{q}(\hat{\mathbf{u}})}{\partial r}, \ t_{n+0.5} < t \le t_{n+1}, \ \hat{\mathbf{u}}|_{t=t_{n+0.5}} = \bar{\mathbf{u}}^{n+0.5}. \tag{157}$$

Let Δr , Δz denote the spatial increments and let $r_i = i\Delta r$, $z_j = j\Delta z$, $i, j = 0, 1, 2, \ldots$. The Runge-Kutta scheme, for Eq. (155) can be written in the form:

$$\hat{\mathbf{v}}_{i,j}^{n+0.75} = \hat{\mathbf{v}}_{i,j}^{n+0.5} - \frac{\Delta t}{2r_i} \mathbf{q} \left(\hat{\mathbf{v}}_{i,j}^{n+0.5} \right),$$

$$\hat{\mathbf{v}}_{i,j}^{n+1} = \hat{\mathbf{v}}_{i,j}^{n+0.5} - \frac{\Delta t}{r_i} \mathbf{q} \left(\hat{\mathbf{v}}_{i,j}^{n+0.75} \right), \quad i = 1, 2, \dots, \quad j = 0, 1, 2, \dots,$$
(158)

where $\hat{\mathbf{v}}_{i,j}^{n+\theta} = \left\{ (\hat{\rho})_{i,j}^{n+\theta}, (\hat{\rho}\hat{v}_r)_{i,j}^{n+\theta}, (\hat{\rho}\hat{v}_z)_{i,j}^{n+\theta}, (\hat{\rho}\hat{E})_{i,j}^{n+\theta} \right\}^T$, $\theta = 0.5, 0.75, 1$. If i = 0, then, in view of (157), the Runge-Kutta scheme is the following

$$\hat{\mathbf{v}}_{0,j}^{n+0.75} = \hat{\mathbf{v}}_{0,j}^{n+0.5} - \frac{\Delta t}{2} \left. \frac{\partial \mathbf{q}}{\partial r} \right|_{r=0}^{n+0.5}, \ \hat{\mathbf{v}}_{0,j}^{n+1} = \hat{\mathbf{v}}_{0,j}^{n+0.5} - \Delta t \left. \frac{\partial \mathbf{q}}{\partial r} \right|_{r=0}^{n+0.75}, \ j \ge 0. \ (159)$$

Notice, we are using, for the sake of convenience, the same notation for the components of, in general, different vectors $\hat{\mathbf{u}}$ and $\hat{\mathbf{v}}$. To derive the scheme at i=0, we will use the reflection concept [2]. We introduce a node i=-1, i.e.,

 $r_{-1} = -\Delta r$. At this node, in view of the reflection concept, we have: $\hat{\mathbf{v}}_{-1,j}^{n+\theta} = \left\{ (\hat{\rho})_{1,j}^{n+\theta}, -(\hat{\rho}\hat{v}_r)_{1,j}^{n+\theta}, (\hat{\rho}\hat{v}_z)_{1,j}^{n+\theta}, (\hat{\rho}\hat{E})_{1,j}^{n+\theta} \right\}^T$, $\theta = 0.5, 0.75, j = 0, 1, 2, \dots$. Then, using central differencing, we obtain, by virtue of (159), the following difference scheme for Eq. (157):

$$\hat{\mathbf{v}}_{0,j}^{n+0.75} = \hat{\mathbf{v}}_{0,j}^{n+0.5} - \frac{\Delta t}{4\Delta r} \left[\mathbf{q} \left(\hat{\mathbf{v}}_{1,j}^{n+0.5} \right) - \mathbf{q} \left(\hat{\mathbf{v}}_{-1,j}^{n+0.5} \right) \right],$$

$$\hat{\mathbf{v}}_{0,j}^{n+1} = \hat{\mathbf{v}}_{0,j}^{n+0.5} - \frac{\Delta t}{2\Delta r} \left[\mathbf{q} \left(\hat{\mathbf{v}}_{1,j}^{n+0.75} \right) - \mathbf{q} \left(\hat{\mathbf{v}}_{-1,j}^{n+0.75} \right) \right], \quad j = 0, 1, 2, \dots$$
(160)

The equations (136)-(139) are integrated up to t=1 with $\sigma \equiv Z_0/R_0=0.1$. It is assumed for the schemes COS2, (59), and MAC2, (64), that the monotonicity parameter $\aleph=0.2$ and the parameter $\varkappa=1$. In view of the stability conditions, (67) we take Cr=0.95 for the scheme COS2. Following to [23], we take Cr=0.475 for the scheme MAC2, in contrast to the sufficient (but not necessary) stability conditions (86). It is assumed, in the case of scheme COS2, that the spatial increments are the following: $\Delta r=10^{-3}$, $\Delta z=5\times10^{-4}$ if $0 < t \le 0.1$; $\Delta r=2\times10^{-3}$, $\Delta z=10^{-3}$ if $0.1 < t \le 0.4$; $\Delta r=2.0\times10^{-3}$, $\Delta z=2.0\times10^{-3}$, if $0.4 < t \le 1$. Taking into account the fact that in the case of the staggered scheme COS2 the space increment is, actually, two times less than in the case of the nonstaggered scheme MAC2, we assume for the scheme MAC2: $\Delta r=5\times10^{-4}$, $\Delta z=2.5\times10^{-4}$ if $0 < t \le 0.1$; $\Delta r=10^{-3}$, $\Delta z=5\times10^{-4}$ if $0.1 < t \le 0.4$; $\Delta r=10^{-3}$, $\Delta z=10^{-3}$ if $0.4 < t \le 1$. The results of simulations as well as the analytical solution are depicted in Figures 7, 8, 9, 10.

We observe (Figures 7, 8, 9, 10) that the numerical results obtained via the schemes COS2 and MAC2 are practically coincide. We also observe that the relatively large deviations between numerical and analytical solutions are occurred in the area of the relatively large second derivatives (i.e. their absolute values) of the primitive variables over the spatial coordinates. Hence, these relatively large deviations may be occurred in the area of small values of the variables in the vicinity of the front (see Figures 7, 8, 9, 10) as well as in the area of relatively large values as well (see Figure 8).

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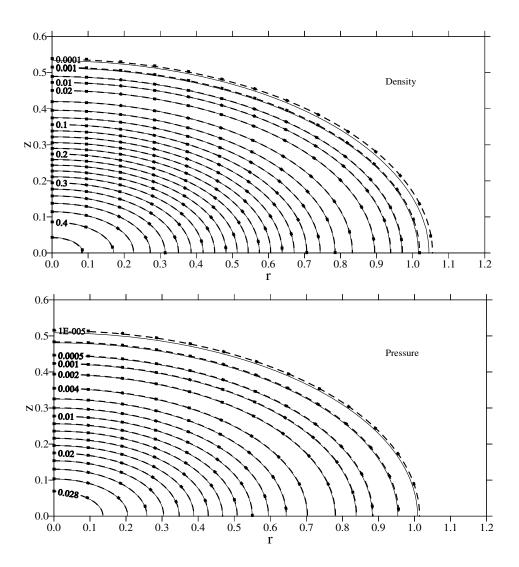


Figure 7: Anisimov's problem, density and pressure distribution. COS2 (Cr = 0.95) and MAC2 (Cr = 0.475) schemes versus the analytical solution. $\sigma \equiv Z_0/R_0 = 0.1$, time t = 0.4. The monotonicity parameter $\aleph = 0.2$ and the parameter $\varkappa = 1$. Dashed lines: numerical solution (COS2); Square dotted lines: numerical solution (MAC2); Solid lines: analytical solution.

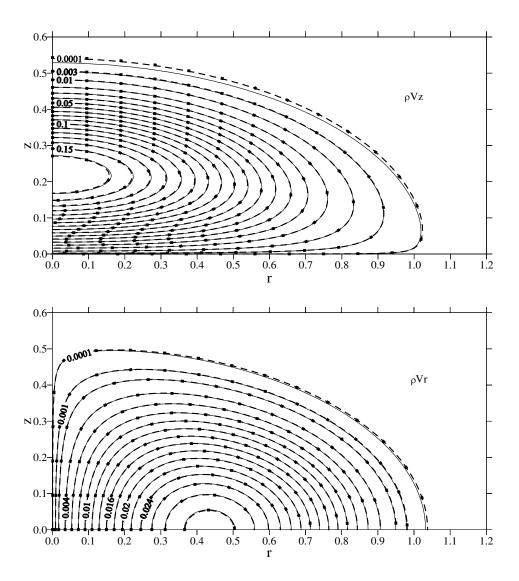


Figure 8: Anisimov's problem, momenta $(\rho V_z \text{ and } \rho V_r)$ distribution. COS2 (Cr=0.95) and MAC2 (Cr=0.475) schemes versus the analytical solution. $\sigma \equiv Z_0/R_0 = 0.1$, time t=0.4. The monotonicity parameter $\aleph=0.2$ and the parameter $\varkappa=1$. Dashed lines: numerical solution (COS2); Square dotted lines: numerical solution (MAC2); Solid lines: analytical solution.

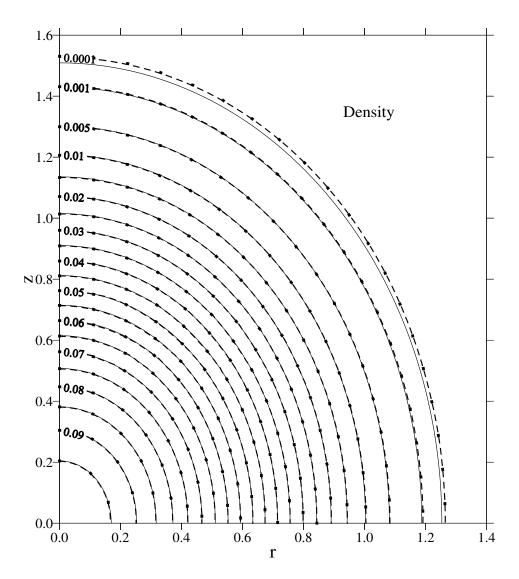


Figure 9: Anisimov's problem, density distribution. COS2 (Cr = 0.95) and MAC2 (Cr = 0.475) schemes versus the analytical solution. $\sigma \equiv Z_0/R_0 = 0.1$, time t = 1. The monotonicity parameter $\aleph = 0.2$ and the parameter $\varkappa = 1$. Dashed lines: numerical solution (COS2); Square dotted lines: numerical solution (MAC2); Solid lines: analytical solution.

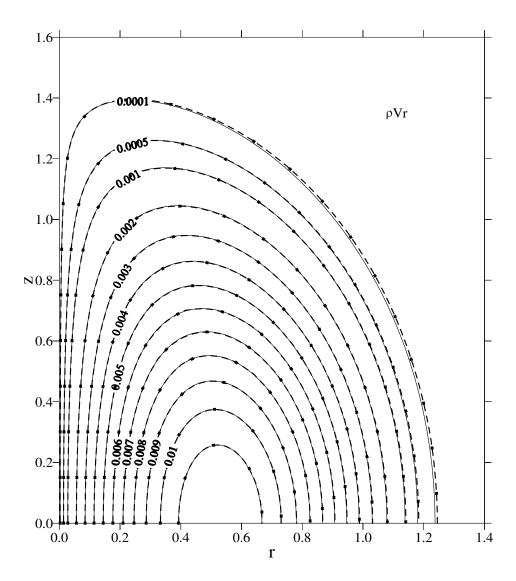


Figure 10: Anisimov's problem, momentum, ρV_r , distribution. COS2 (Cr=0.95) and MAC2 (Cr=0.475) schemes versus the analytical solution. $\sigma\equiv Z_0/R_0=0.1$, time t=1. The monotonicity parameter $\aleph=0.2$ and the parameter $\varkappa=1$. Dashed lines: numerical solution (COS2); Square dotted lines: numerical solution (MAC2); Solid lines: analytical solution.

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